

Estimating Failure Probabilities

Holger Drees* and Laurens de Haan†

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Abstract

In risk management often the probability must be estimated that a random vector falls into an extreme failure set. In the framework of bivariate extreme value theory, we construct an estimator for such failure probabilities and analyze its asymptotic properties under natural conditions. It turns out that the estimation error is mainly determined by the accuracy of the statistical analysis of the marginal distributions. Moreover, we establish confidence intervals and briefly discuss generalizations to higher dimensions and issues arising in practical applications as well.

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1 Introduction

Suppose an insurance company has contracts in two related lines of business with all customers of an insurance portfolio (e.g., fire insurance and business interruption insurance for industrial customers). On top of quota reinsurances for both lines of business (possibly with different quotas) the remaining total loss from each incidence is covered by an excess of loss reinsurance (CAT-XL) that pays for the part of the total loss which exceeds a given high retention level R . If X and Y denote the original losses from a fire in both lines of business and $1 - \alpha_X$ and $1 - \alpha_Y$ the corresponding quotas, then a claim occurs in the XL-reinsurance if $\alpha_X X + \alpha_Y Y$ exceeds R . For the purpose of risk management the reinsurer might be interested in the probability that the insurance company will file a claim in case of a fire. If the retention level is high, then such the claim probability cannot be estimated using simple empirical estimates, because in the past the retention has rarely (or never) been exceeded.

In this paper a more general setting is considered. We are interested in estimating the probability that a pair of random variables (X, Y) will take on a value in some given “extreme” set. Similar problems arise naturally in many fields. For example, a coastal dike may fail if the vector build from the still water level and the wave heights lie in a certain failure set D (cf. Coles and Tawn,

*University of Hamburg, Department of Mathematics, SPST, Bundesstr. 55, 20146 Hamburg, Germany; email: holger.drees@uni-hamburg.de

†Erasmus University Rotterdam, Department of Economics, P.O. Box 1738, 3000 DR Rotterdam, The Netherlands; email: ldehaan@ese.eur.nl

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1994, Bruun and Tawn, 1998, and de Haan and de Ronde, 1998). A financial option (like a down-and-out-put) may become worthless if the price vector of underlyings enters such a “failure set”. Finally, (part of) the principal of a catastrophe bond gets lost for the investors if a vector of triggers becomes too extreme.

As there are insufficiently many observations available in the extreme failure set D to use standard statistical methods, extreme value theory is needed to estimate the failure probability $P\{(X, Y) \in D\}$. The basic idea of multivariate extreme value theory is to assume that the suitably standardized componentwise maxima of the observed random variables converge to a non-degenerate limit distribution. It can be shown that this assumption is equivalent to the convergence of suitably standardized quantile functions of both marginal distributions and a condition on the dependence structure in extreme regions.

To be more precise, denote the marginal distribution functions of X and Y by F_1 and F_2 , respectively, and let $U_i(t) := F_i^{\leftarrow}(1 - 1/t)$ with $H^{\leftarrow}(s) := \inf\{x \in \mathbb{R} \mid H(x) \geq s\}$ denoting the generalized inverse of an increasing function H . We assume that there exist real constants γ_i , positive functions a_i and real functions b_i such that for $x > 0$ and $i \in \{1, 2\}$

$$\lim_{t \rightarrow \infty} \frac{U_i(tx) - b_i(t)}{a_i(t)} = \frac{x^{\gamma_i} - 1}{\gamma_i}. \quad (1.1)$$

For $\gamma_i = 0$ read the right-hand side as $\log x$. Note that the right-hand side is the U -function of the generalized Pareto distribution (GPD) with distribution function $1 - (1 + \gamma_i x)^{-1/\gamma_i}$ for $1 + \gamma_i x > 0$, that is to be interpreted as the standard exponential distribution function for $\gamma_i = 0$. The parameter γ_i is the so-called extreme value index of the i th marginal. If it is positive, then the support of F_i is unbounded from above and $1 - F_i(t)$ roughly decays like the power function with exponent $1/\gamma_i$, while for $\gamma_i < 0$ the right endpoint $x_i^* := F_i^{\leftarrow}(1)$ of the support is finite and $1 - F_i(x)$ roughly behaves like a multiple of $(x_i^* - x)^{-1/\gamma_i}$ as $x \uparrow x_i^*$.

The aforementioned extremal dependence condition can be given in terms of the standardized random variables $1 - F_1(X)$ and $1 - F_2(Y)$, that are uniformly distributed on $[0, 1]$ if the marginal distributions are continuous. More precisely, we assume the existence of a measure ν such that for ν -continuous Borel sets $B \subset [0, \infty)^2$ bounded away from the origin

$$\lim_{t \rightarrow \infty} tP\{(X, Y) \in U(tB)\} = \nu(B). \quad (1.2)$$

Here and in what follows, for functions h_1, h_2 which are defined on subsets of the reals, we define a function h on a subset of \mathbb{R}^2 by $h(x_1, x_2) := (h_1(x_1), h_2(x_2))$. The so-called exponent measure ν describes the asymptotic dependence structure between extreme observations X and Y . Its homogeneity property

$$\nu(tB) = t^{-1}\nu(B), \quad (1.3)$$

which holds for all Borel sets $B \subset [0, \infty)^2$ and all $t > 0$, will be pivotal for the construction of our estimator of the failure probability. (Seen from a different angle, we assume an approximate scaling law for the joint distribution of $U^{\leftarrow}(X, Y)$; cf. Anderson (1994).) Further details about the extreme value assumptions can be found in de Haan and Ferreira (2006), Sections 1.2 and 6.1, or Beirlant et al. (2004), Chapters 2 and 8.

We are interested in the situation that at most a few observations lie in the extreme failure set D which implies that in our mathematical framework the failure set $D = D_n$ must depend on the sample size n such that the failure probability

$$p_n := P\{(X, Y) \in D_n\}$$

tends to 0. To motivate an estimator of p_n based on independent copies (X_i, Y_i) , $1 \leq i \leq n$, of (X, Y) first note that from (1.2) we obtain the approximation

$$\frac{n}{k} P\left\{\frac{k}{n} U^\leftarrow(X, Y) \in B\right\} \approx \nu(B) \quad (1.4)$$

for any sequence $k = k_n \rightarrow \infty$ such that $k/n \rightarrow 0$. To estimate p_n using this approximation, we must replace U^\leftarrow and ν with suitable estimators.

According to (1.1), we may approximate $U_i((n/k)x)$ for sufficiently large n by

$$T_{n,i}(x) := a_i(n/k) \frac{x^{\gamma_i} - 1}{\gamma_i} + b_i(n/k) \quad (1.5)$$

and estimate it by

$$\hat{T}_{n,i}(x) := \hat{a}_i(n/k) \frac{x^{\hat{\gamma}_i} - 1}{\hat{\gamma}_i} + \hat{b}_i(n/k), \quad (1.6)$$

where $\hat{a}_i(n/k)$, $\hat{b}_i(n/k)$ and $\hat{\gamma}_i$ are suitable estimators for $a_i(n/k)$, $b_i(n/k)$ and γ_i , respectively. (Later on, the joint asymptotic normality of these estimators is needed; see e.g. de Haan and Ferreira, 2006, Chapters 3 and 4, or Beirlant et al., 2004, Section 5.6. for examples.) Likewise, the generalized inverse functions $(k/n)U_i^\leftarrow(x)$ can be estimated by

$$\hat{T}_{n,i}^\leftarrow(x) := \left(1 + \hat{\gamma}_i \frac{x - \hat{b}_i(n/k)}{\hat{a}_i(n/k)}\right)^{1/\hat{\gamma}_i}. \quad (1.7)$$

Here and in the sequel, $(1 + \gamma y)^{1/\gamma}$ is defined as e^y if $\gamma = 0$. For $1 + \gamma y < 0$ (or $1 + \gamma y = 0$ and $\gamma < 0$) the term $(1 + \gamma y)^{1/\gamma}$ is not well defined. If γ is positive and $y < -1/\gamma$, then it may be interpreted as 0, while for $\gamma < 0$ and $y > -1/\gamma$ it may be defined to be ∞ . However, we will see that the precise definition of $(1 + \gamma y)^{1/\gamma}$ for very small and for negative values of $1 + \gamma y$ is not important in the present setting (provided it is taken to be a non-decreasing function of y), because the sets on which $\hat{T}_{n,i}^\leftarrow$, $i \in \{1, 2\}$, are not well defined are asymptotically negligible.

If, in (1.4), we substitute $\hat{T}_n^\leftarrow(x_1, x_2) := (\hat{T}_{n,1}^\leftarrow(x_1), \hat{T}_{n,2}^\leftarrow(x_2))$ for the marginal transformation $(k/n)U^\leftarrow$ and replace the probability in the left-hand side of (1.4) by its empirical counterpart, we arrive at the following estimator of ν

$$\hat{\nu}_n(B) := \frac{1}{k} \sum_{i=1}^n \varepsilon_{\hat{T}_n^\leftarrow(X_i, Y_i)}(B), \quad (1.8)$$

with ε_x denoting the Dirac measure with mass 1 at x .

Now, again interpreting convergence (1.2) (for $t = e_n$) as an approximation, we may estimate the failure probability as follows:

$$\begin{aligned} p_n &= P\{(X, Y) \in D_n\} \\ &= P\{(X, Y) \in U(e_n \cdot e_n^{-1} U^\leftarrow(D_n))\} \\ &\approx \frac{1}{e_n} \nu(e_n^{-1} U^\leftarrow(D_n)) \end{aligned} \quad (1.9)$$

$$\begin{aligned} &\approx \frac{1}{e_n} \nu\left(\frac{n}{k e_n} \hat{T}_n^\leftarrow(D_n)\right) \\ &\approx \frac{1}{e_n} \hat{\nu}_n\left(\frac{n}{k e_n} \hat{T}_n^\leftarrow(D_n)\right) \\ &=: \hat{p}_n. \end{aligned} \quad (1.10)$$

The basic idea of this estimator is to blow up the failure set, after a standardization of the marginals, such that it contains sufficiently many observations to allow the estimation of its probability by an empirical probability. Here the constants k and e_n , that control by how much the transformed failure set is inflated, can be chosen by the statistician. This must be done appropriately to balance two contrary effects. On the one hand, ke_n must not be too small, such that the inflated standardized failure set $n/(ke_n)\hat{T}_n^{\leftarrow}(D_n)$ contains sufficiently many marginally transformed observations $\hat{T}_n^{\leftarrow}(X_i, Y_i)$, and thus the empirical probability (1.8) is an accurate estimate of its expectation. On the other hand, the set $e_n^{-1}U^{\leftarrow}(D_n)$ must be sufficiently extreme to justify approximation (1.9). In Section 3 we discuss a heuristic tool to ensure this balance.

The main goal of the present paper is to establish the asymptotic normality of the estimator \hat{p}_n under conditions on the underlying distribution and the failure set which are easy to interpret and relatively simple to verify. In what follows we outline how to achieve this objective.

Recall that, in our asymptotic framework, the failure set D_n must become more extreme in the sense that it moves in the north-east direction as the sample size n increases to ensure that it contains at most a few observations. To make both coordinates comparable, we standardize the marginals using U^{\leftarrow} and assume that $U^{\leftarrow}(D_n)$ is essentially an increasing multiple of a fixed set S . That way we ensure that none of the coordinates dominates the other. More precisely, we assume that for different sample sizes the failure sets are obtained from one fixed set $S \subset [0, \infty)^2$ in that there exist constants $d_n > 0$ tending to ∞ such that

$$D_n = U(d_n S) \cap \mathbb{R}^2 = \{(U_1(d_n x), U_2(d_n y)) \mid (x, y) \in S\} \cap \mathbb{R}^2. \quad (1.11)$$

Note that from the analog to (1.9) where e_n is replaced with d_n one obtains $d_n \approx \nu(S)/p_n$ (see Lemma 4.9 for a precise proof of the assertion $p_n d_n \rightarrow \nu(S)$). Hence the model constants d_n determine at which rate the failure probabilities tend to 0.

The crucial idea in the analysis of the asymptotic behavior of \hat{p}_n is to approximate the estimator by the empirical measure of a *random transformation* $H_n(S)$ of the set S (with H_n defined in (1.12) below) under the following analog to $\hat{\nu}_n$ (defined in (1.8)) with the fitted GPDs replaced by the “true” ones:

$$\nu_n(B) := \frac{1}{k} \sum_{i=1}^n \varepsilon_{T_n^{\leftarrow}(X_i, Y_i)}(B).$$

Since the GPD approximation of the marginals is accurate only in the upper tail (and to avoid the aforementioned problem with the definition of T_n^{\leftarrow}), we must first show that asymptotically it does not matter if we replace S with a suitably defined subset S_n^* that is bounded away from the coordinate axes. For this set, we may use the approximation

$$\hat{p}_n \approx \frac{1}{e_n} \nu_n\left(\frac{d_n}{e_n} H_n(S_n^*)\right).$$

where the random transformation H_n of the marginals is defined by

$$H_n(x) := \frac{e_n}{d_n} T_n^{\leftarrow} \circ \hat{T}_n \circ (\hat{T}_n^{(c)})^{\leftarrow} \circ U(d_n x) \quad (1.12)$$

with

$$c = c_n := \frac{k}{n} e_n \quad (1.13)$$

and

$$\hat{T}_n^{(c)}(x, y) = \hat{T}_n(c_n x, c_n y). \quad (1.14)$$

Check that by (1.1) one has $H_n(x) \approx (e_n/d_n)(T_n^{(c)})^{\leftarrow} \circ U(d_n x) \approx (e_n/d_n)(T_n^{(c)})^{\leftarrow} \circ T_n((k/n)d_n x) \approx x$ (cf. Lemma 4.1). Using the homogeneity of ν , we will thus break the estimation error into several parts as follows:

$$\begin{aligned} \hat{p}_n - p_n &= \hat{p}_n - \frac{1}{e_n} \nu_n\left(\frac{d_n}{e_n} H_n(S_n^*)\right) \\ &\quad + \frac{1}{e_n} (\nu_n(B) - E\nu_n(B))|_{B=(d_n/e_n)H_n(S_n^*)} \\ &\quad + \frac{1}{e_n} (E\nu_n(B) - \nu(B))|_{B=(d_n/e_n)H_n(S_n^*)} \\ &\quad + \frac{1}{d_n} (\nu(H_n(S_n^*)) - \nu(S_n^*)) \\ &\quad + \frac{1}{d_n} (\nu(S_n^*) - \nu(S)) \\ &\quad + \nu(d_n S) - p_n \\ &=: I + II + III + IV + V + VI. \end{aligned} \quad (1.15)$$

It will turn out that, under suitable conditions, part IV dominates all the other terms. Its asymptotic behavior is largely determined by the asymptotics of the marginal estimators if ν is sufficiently smooth.

Under very weak conditions on the set S , we will show that the terms I and V are negligible, if S_n^* is defined suitably. If d_n/e_n is bounded and bounded away from 0, then using methods from empirical process theory the second term can be shown to be asymptotically negligible. Part VI is a bias term which is negligible if d_n is sufficiently large (depending on the rate of convergence in (1.2)). Similarly, the term III, that equals $((n/k)P\{T_n^{\leftarrow}(X, Y) \in B\} - \nu(B))/d_n$ for $B = H_n(S_n^*)$, describes a bias term which is asymptotically negligible if both the approximation (1.4) and the marginal approximation $U((n/k)B) \approx T_n(B)$ are sufficiently accurate.

An estimator related to \hat{p}_n has been suggested and analyzed by de Haan and Sinha (1999) in a much more restrictive framework. In particular, it has been assumed that the same number k_n of largest order statistics is used for both marginal fits, which is inefficient if the GPD approximation (cf. (2.1) below) is less accurate for one of the marginal distributions. Likewise, the flexibility of the estimator is increased in the present paper by allowing that the blow-up factor e_n deviates from the unknown model constant d_n , while de Haan and Sinha (1999) used a consistent estimator of d_n , which was made identifiable by fixing some point on the boundary of S .

Moreover, the shape of the failure set considered by de Haan and Sinha is restricted. E.g. the case $q(\infty) = 0$ (in our notation; cf. condition (Q2) below) is ruled out by condition (2.9) of that paper. The model assumption

$$D_n := \{(s, t) \mid f(s/x_n, t/y_n) \geq 1\}$$

for some function f and sequences of normalizing constants x_n and y_n seems quite restrictive and unnatural, because it allows the failure set to tend towards the “north-east” only by a linear scaling of both marginals. This parametrization does not fit well to extreme value theory if the extreme value indices are not positive, which is usually the case in environmetrics, arguably the most important field of application of our theory. Even more troublesome is the fact that by

assumption (1.5) of de Haan and Sinha (1999) the failure set is described in terms of the number k_n of largest order statistics that is picked by the statistician. Hence the model parametrization depends on the statistical procedure used to analyze the model, which makes it extremely difficult to interpret.

Finally, while the influence of each marginal transformation is clearly separated in the description of the limiting distribution in our main Theorem 2.1, in Theorem 4.1 of de Haan and Sinha (1999) the marginal parameters are seemingly intermingled. Therefore, the generalization of the present results to higher dimension is much more straightforward than those of de Haan and Sinha (see the discussion at the end of Section 2).

An alternative to our genuinely multivariate estimator can be constructed by the so-called structural variable approach if the failure set is of the form $D_n = \{(s, t) \mid h(s, t) \geq t_n\}$ for some known function h and threshold t_n . Then one may apply techniques from univariate extreme value theory to the pseudo-observations $h(X_i, Y_i)$, $1 \leq i \leq n$ (cf. Coles (2001), §8.2.4 and page 156, or Bruun and Tawn (1998)). However, even for this class of failure sets, an analysis of the dependence structure between the two components of the observed vectors is of independent interest, and it seems more natural to use the same approach for model fitting and for the estimation of quantities like failure probabilities. Moreover, often one wants to estimate the failure probability for several different sets (e.g., to find the cheapest construction to ensure a certain level of safety); in this case it is both more efficient and more natural to use estimators in a unified framework as considered in the present paper.

In the multivariate approach, Coles and Tawn (1994) and Bruun and Tawn (1998) used parametric models for the dependence structure in the closely related problem to estimate a parameter defining a failure set such that the corresponding failure probability equals a given value. However, usually there is no physical reason for such parametric models. By using them nevertheless one trades a modeling error, that is difficult to assess, for an estimation error, which can be quantified at least asymptotically (see Theorem 2.1 below). Indeed, Davison (1994) suggested in the discussion to Coles and Tawn (1994) that for sufficiently large sample sizes a nonparametric approach may be advisable.

Finally, we would like to mention that the assumptions used in the present paper rule out that the exponent measure ν concentrates on the coordinate axes (i.e., here X and Y are assumed asymptotically dependent in the sense of multivariate extreme value theory). In the case of asymptotic independent coordinates X and Y , consistency of an analogous estimator for the failure probability was proved by Draisma et al. (2004), while its asymptotic normality was established by Müller (2008).

The paper is organized as follows: In Section 2 we first introduce and discuss in detail the framework in which we then prove asymptotic normality of our estimator of the failure probability. Moreover, we propose a consistent estimator of the limiting variance and derive an asymptotic confidence interval. In Section 3 we apply the theory to the motivating example given at the beginning. In this context, we also discuss the roles of k and e_n and propose a heuristic approach for choosing those numbers. All proofs are collected in Section 4.

2 Main results

We will make the following assumptions about the marginal distributions and the estimators of the marginal parameters:

(M1) There exist constants $x_i^0 < F_i^{\leftarrow}(1)$ such that F_i is continuous and strictly increasing on $[x_i^0, F_i^{\leftarrow}(1)] \cap \mathbb{R}$ for $i \in \{1, 2\}$.

(M2) For all $i \in \{1, 2\}$, there exist normalizing functions $a_i > 0$, $b_i \in \mathbb{R}$ and $A_i \neq 0$ and constants $\rho_i < 0$ such that for all $x > 0$

$$\lim_{t \rightarrow \infty} \frac{\frac{U_i(tx) - b_i(t)}{a_i(t)} - \frac{x^{\gamma_i - 1}}{\gamma_i}}{A_i(t)} = \bar{\psi}_{\gamma_i, \rho_i}(x) := \begin{cases} \frac{x^{\gamma_i + \rho_i}}{\gamma_i + \rho_i}, & \gamma_i + \rho_i \neq 0 \\ \log x, & \gamma_i + \rho_i = 0 \end{cases}$$

(M3)

$$k^{1/2} \left(\frac{\hat{a}_i(n/k)}{a_i(n/k)} - 1, \frac{\hat{b}_i(n/k) - b_i(n/k)}{a_i(n/k)}, \hat{\gamma}_i - \gamma_i \right)_{1 \leq i \leq 2} \longrightarrow (\alpha_i, \beta_i, \Gamma_i)_{1 \leq i \leq 2}$$

weakly.

Condition (M1) is not essential, but it is assumed to simplify the proofs and the formulation of some technical results (cf. de Haan and Ferreira, 2006, Theorem B.3.13).

(M2) is the usual second order condition with the additional restriction that the second order parameters ρ_i are negative. Again, one may drop the latter assumption at the cost of additional technical complications. According to Corollary 2.3.7 of de Haan and Ferreira (2006) we may and will assume that the normalizing constants are chosen such that the following uniform version holds: For all $\varepsilon, \delta > 0$ there exists t_0 such that

$$\left| \frac{\frac{U_i(tx) - b_i(t)}{a_i(t)} - \frac{x^{\gamma_i - 1}}{\gamma_i}}{A_i(t)} - \bar{\psi}_{\gamma_i, \rho_i}(x) \right| \leq \delta x^{\gamma_i + \rho_i} \max(x^\varepsilon, x^{-\varepsilon}) =: \delta x^{\gamma_i + \rho_i \pm \varepsilon} \quad (2.1)$$

provided $t, tx > t_0$. In fact, the main results hold under the following weaker assumption:

$$\left| \frac{U_i(tx) - b_i(t)}{a_i(t)} - \frac{x^{\gamma_i - 1}}{\gamma_i} \right| = O(A_i(t) x^{\gamma_i + \rho_i \pm \varepsilon}) \quad (2.2)$$

as $t \rightarrow \infty$ uniformly for $x \geq t_0/t$. Note that, under condition (M2), A_i is regularly varying with index ρ_i .

Condition (M3) gives a lower bound on the rate at which the marginal estimators converge. Notice that some of the limiting random variables may be equal to 0 almost surely. In particular, this will usually be the case, if the i th marginal estimators use k_i largest order statistics and $k = o(k_i)$. However, typically at least some of the limiting random variables are non-degenerate and jointly normal distributed. In the sequel, we will choose versions such that the convergence in (M3) holds in probability.

The failure set D_n has to satisfy the following conditions.

(Q1) There exists a set

$$S = \{(x, y) \in [0, \infty)^2 \mid y \geq q(x) \forall x \in [0, \infty)\} \subset [0, \infty)^2$$

and constants $d_n > 0$ tending to ∞ such that

$$D_n = U(d_n S) \cap \mathbb{R}^2 = \{(U_1(d_n x), U_2(d_n y)) \mid (x, y) \in S\} \cap \mathbb{R}^2.$$

Here $q : [0, \infty) \rightarrow [0, \infty]$ is some function that is monotonically decreasing and continuous from the right with $q(0) > 0$.

(Q2)

$$\begin{aligned} x^{(1-\gamma_1)/2} |\log x| &= O(q(x)) \quad \text{as } x \downarrow x_l := \inf\{x \geq 0 \mid q(x) < \infty\} \\ y^{(1-\gamma_2)/2} |\log y| &= O(q^\leftarrow(y)) \quad \text{as } y \downarrow q(\infty) := \lim_{x \rightarrow \infty} q(x). \end{aligned}$$

In particular, condition (Q1) ensures that one may define the generalized (right-continuous) inverse function (of a decreasing function) in the usual way:

$$q^\leftarrow(v) := \inf\{x > 0 \mid q(x) \leq v\}$$

with the convention $\inf \emptyset = \infty$. The conditions (Q2) are always fulfilled if $\gamma_1 \leq 1$ or $x_l > 0$, resp., if $\gamma_2 \leq 1$ or $q(\infty) > 0$.

Moreover, we need some conditions on the extremal dependence between X and Y . Recall that asymptotically the extremal dependence is described by the exponent measure ν defined in (1.2). In view of (1.1), one may replace the standardization by U with a standardization using T_n . To bound the bias terms III and VI in (1.15), we must specify the rate of the resulting convergence towards ν :

(D1) There exist an exponent measure ν on $[0, \infty)^2$ and a function $A_0(t) > 0$ converging to 0 as t tends to ∞ such that

$$t_n P\left\{\left(\left(1 + \gamma_1 \frac{X - b_1(t_n)}{a_1(t_n)}\right)^{1/\gamma_1}, \left(1 + \gamma_2 \frac{Y - b_2(t_n)}{a_2(t_n)}\right)^{1/\gamma_2}\right) \in B\right\} - \nu(B) = O(A_0(t_n))$$

uniformly for all sets $B \in \mathcal{B}_{t_n, M}$ for $t_n = n/k$ and for $t_n = d_n$ and arbitrary $M > 0$.

Here, $\mathcal{B}_{t_n, M}$ consists of all sets of the form $\{(\tilde{H}_{\vartheta_i, \chi_i, \xi_i}^{(n, i)}(x_i))_{i \in \{1, 2\}} \mid (x_1, x_2) \in C\}$ with $C = S \cap [u, \infty) \times [v, \infty)$ or $C = [x_l, u) \times [q(u-), \infty)$ or $C = [q^\leftarrow(v), \infty) \times [q(\infty), v)$ for some $u, v > 0$ and some $\vartheta_i, \chi_i, \xi_i \in [-M, M]$ if $t_n = n/k$, and $\mathcal{B}_{t_n, M}$ comprises all sets of the form $\{((1 + \gamma_i(U_i(d_n x_i) - b_i(d_n))/a_i(d_n))^{1/\gamma_i})_{i \in \{1, 2\}} \mid (x_1, x_2) \in C\}$ with $C = [x_l, u) \times [q(u-), \infty)$ or $C = [q^\leftarrow(v), \infty) \times [q(\infty), v)$ for some $u, v > 0$ if $t_n = d_n$. Here, for $i \in \{1, 2\}$,

$$\tilde{H}_{\vartheta_i, \chi_i, \xi_i}^{(n, i)}(x) := \left[1 + \gamma_i \left(\frac{c_n^{-(\gamma_i - k^{-1/2}\vartheta_i)} - 1}{\gamma_i - k^{-1/2}\vartheta_i} (1 + k^{-1/2}\xi_i) + c_n^{-(\gamma_i + k^{-1/2}\chi_i)} \frac{U_i(d_n x) - b_i(n/k)}{a_i(n/k)}\right)\right]^{1/\gamma_i}. \quad (2.3)$$

(D2) The exponent measure has a strictly positive, continuous Lebesgue density η .

Alternatively, the dependence may be described by the pertaining spectral measure Φ on $[0, \pi/2]$ defined by

$$\Phi([0, \vartheta]) = \nu\left\{(x, y) \in [0, \infty)^2 \mid x^2 + y^2 > 1, \arctan \frac{y}{x} \leq \vartheta\right\}, \quad \vartheta \in [0, \pi/2].$$

Condition (D2) is equivalent to the assumption that Φ has a continuous Lebesgue density φ that is strictly positive on $[0, \pi/2]$. In particular, this assumption rules out that X and Y are asymptotically independent (in the sense of multivariate extreme value theory), because then the spectral measure is concentrated on $\{0, \pi/2\}$. The relationship between η and φ is given by

$$\eta(x, y) = (x^2 + y^2)^{-3/2} \varphi\left(\arctan \frac{y}{x}\right), \quad x, y > 0. \quad (2.4)$$

In (D1) the rectangles can also be replaced with the subsets $S \cap ((0, u) \times (0, \infty))$, resp. $S \cap ((0, \infty) \times (0, v))$. It is easy to see that condition (D1) is met if (X, Y) has a density f which satisfies the following approximation

$$\sup_{(x, y) \in (0, \infty)^2, x \vee y \geq 1} \frac{1}{w(x, y)} \left| t a_1(t) a_2(t) x^{\gamma_1 - 1} y^{\gamma_2 - 1} f\left(a_1(t) \frac{x^{\gamma_1} - 1}{\gamma_1} + b_1(t), a_2(t) \frac{y^{\gamma_2} - 1}{\gamma_2} + b_2(t)\right) - \eta(x, y) \right| = O(A_0(t))$$

for some weight function w which is Lebesgue-integrable on $\{(x, y) \in (0, \infty)^2, x \vee y \geq 1\}$. This sufficient condition applies e.g. to the bivariate Cauchy distribution restricted to $(0, \infty)^2$ and to densities of the form $f(x, y) = 1/(1 + x^\alpha + y^\beta)$ with $\alpha, \beta > 1$ such that $\beta > \alpha/(\alpha - 1)$.

Finally, we impose the following conditions on the sequences d_n, e_n and $k = k_n$:

(S1) $k \rightarrow \infty$, $n = O(e_n)$ (so that $k = O(c_n)$ with $c_n = e_n k/n \rightarrow \infty$), $d_n \asymp e_n$ (i.e. $0 < \liminf d_n/e_n \leq \limsup_{n \rightarrow \infty} d_n/e_n < \infty$), and $w_n(\gamma_i) = o(k^{1/2})$ for $i \in \{1, 2\}$ with

$$w_n(\gamma_i) := \begin{cases} \log c_n, & \gamma_i > 0 \\ \frac{1}{2} \log^2 c_n, & \gamma_i = 0 \\ (d_n k/n)^{-\gamma_i}, & \gamma_i < 0. \end{cases}$$

(S2) $A_i(n/k) = o(k^{-1/2} w_n(\gamma_i))$ for $i \in \{1, 2\}$ and $A_0(n/k) = o(k^{-1/2} \max(w_n(\gamma_1), w_n(\gamma_2)))$

(S3) $k^{1/2} = O(c_n \vee c_n^{\gamma_i})$ if $\gamma_i \geq 0$ for $i \in \{1, 2\}$,
 $k^{1/2} = o(c_n^{1-\gamma_1})$ if $\gamma_1 < 0$ and $x_l = 0$, and
 $k^{1/2} = o(c_n^{1-\gamma_2})$ if $\gamma_2 < 0$ and $q(\infty) = 0$.

Recall that d_n is a constant determined by the model, which describes the rate at which the probability p_n to be estimated tends to 0, while e_n is chosen by the statistician such that the inflated failure set contains sufficiently many observations. It seems natural to choose e_n of the same order as d_n , because this way one compensates for the shrinkage of D_n . More precisely, $d_n \asymp e_n$ if and only if the expected number of transformed observations in the inflated transformed failure set is of the same order as k , which can easily be checked in practical applications. To see this, note that by (1.4), (1.3) and (1.11) this expected number equals

$$nP\left\{\hat{T}_n^{\leftarrow}(X, Y) \in \frac{n}{k e_n} \hat{T}_n^{\leftarrow}(D_n)\right\} \approx nP\left\{\frac{k}{n} U^{\leftarrow}(X, Y) \in \frac{1}{e_n} U^{\leftarrow}(D_n)\right\} \approx k \nu\left(\frac{d_n}{e_n} S\right) = k \frac{e_n}{d_n} \nu(S).$$

We would like to emphasize, though, that this condition can be substantially weakened at the price that one needs different conditions for different combinations of signs of γ_1 and γ_2 .

The first condition of (S1) ensures that the expected number of marginally standardized observations in the *inflated* standardized failure region tends to ∞ , whereas the second condition means

that the expected number of observations in the failure region remains bounded as $n \rightarrow \infty$. The last condition of (S1) is needed to ensure consistency of the estimator in the sense that $\hat{p}_n/p_n \rightarrow 1$. Note that it can only be satisfied if $\min(\gamma_1, \gamma_2) > -1/2$. This restriction on the extreme value indices usually arises if one wants to prove asymptotic normality for estimators of tail probabilities; cf., e.g., de Haan and Ferreira (2006), Remark 4.4.3, or Drees et al. (2006), Remark 2.2.

From (S2) it follows that the bias is asymptotically negligible, while (S3) will imply that the part of the set S near the axes (corresponding to observations where one of the coordinates is much larger than the other) does not play an important role asymptotically. Similarly as above, these conditions may also be substantially weakened at the price of much more complicated conditions on the behavior of q depending on γ_1, γ_2 and η .

Under these conditions we establish an asymptotic approximation of the estimator \hat{p}_n .

Theorem 2.1. *If the conditions (M1)–(M3), (D1), (D2), (Q1), (Q2) and (S1)–(S3) are fulfilled, then*

$$\begin{aligned}
& k^{1/2} d_n (\hat{p}_n - p_n) \\
&= w_n(\gamma_1) \begin{cases} -\frac{\Gamma_1}{\gamma_1} \int_{q(\infty)}^{\infty} q^{\leftarrow}(v) \eta(q^{\leftarrow}(v), v) dv, & \gamma_1 > 0 \\ (\frac{\alpha_1}{\gamma_1} - \beta_1 - \frac{\Gamma_1}{\gamma_1}) \int_{q(\infty)}^{\infty} (q^{\leftarrow}(v))^{1-\gamma_1} \eta(q^{\leftarrow}(v), v) dv, & \gamma_1 < 0 \\ -\Gamma_1 \int_{q(\infty)}^{\infty} q^{\leftarrow}(v) \eta(q^{\leftarrow}(v), v) dv, & \gamma_1 = 0 \end{cases} \\
&+ w_n(\gamma_2) \begin{cases} -\frac{\Gamma_2}{\gamma_2} \int_{x_l}^{\infty} q(u) \eta(u, q(u)) du, & \gamma_2 > 0 \\ (\frac{\alpha_2}{\gamma_2} - \beta_2 - \frac{\Gamma_2}{\gamma_2}) \int_{x_l}^{\infty} (q(u))^{1-\gamma_2} \eta(u, q(u)) du, & \gamma_2 < 0 \\ -\Gamma_2 \int_{x_l}^{\infty} q(u) \eta(u, q(u)) du, & \gamma_2 = 0 \end{cases} \\
&+ o_P(w_n(\gamma_1) \vee w_n(\gamma_2))
\end{aligned} \tag{2.5}$$

Since $p_n d_n \rightarrow \nu(S)$, Theorem 2.1 remains true when the left-hand side of (2.5) is replaced with $k^{1/2} \nu(S)(\hat{p}_n/p_n - 1)$.

The weights $w_n(\gamma_1)$ and $w_n(\gamma_2)$ on the right-hand side of (2.5) may be different, and then they converge to ∞ at different rates. More precisely, $w_n(\gamma)$ is a non-increasing function of γ , and it is strictly decreasing on $(-\infty, 0]$. Therefore, the smaller of both marginal extreme value indices γ_1 and γ_2 determines the rate of convergence of \hat{p}_n towards p_n . If at least one of the indices is non-positive and the indices are not equal, then the summand corresponding to the larger index is negligible. (In that case, it may happen that one cannot prove asymptotic normality using Theorem 2.1, because the limiting random variables α_i, β_i and Γ_i pertaining to the smaller extreme value index are equal to 0; cf. the above discussion of condition (M3).)

If both extreme value indices are positive, then both main terms on the right-hand side of (2.5) are of the same order. In that case, $(k^{1/2} d_n / \log c_n)(\hat{p}_n - p_n)$ converge to a limit distribution which typically will be non-degenerate if at least one of the limiting random variables Γ_1 and Γ_2 in (M3) is non-degenerate. If they are jointly normal, then we may derive the asymptotic normality of the estimator for the failure probability p_n .

Theorem 2.1 can be used to construct asymptotic confidence intervals. To this end, it is advisable to reformulate the assertion as a convergence result on $k^{1/2} e_n(\hat{p}_n - p_n)$, because d_n is unknown. Then one needs consistent estimators for the variance of the random variables occurring on the right-hand side of (2.5) which usually are asymptotically normal, and consistent estimators for e_n/d_n times the integral there.

We will outline how to estimate the term $(e_n/d_n) \int_{x_l}^{\infty} q(u)\eta(u, q(u)) du$, that is needed in the case $\gamma \geq 0$. To avoid the estimation of the density η of ν , we approximate the integral by the ν -measure of a shrinking set as follows. Because η is continuous, for small ℓ_n one has

$$\frac{e_n}{d_n} \int_{x_l}^{\infty} q(u)\eta(u, q(u)) du \approx \frac{e_n}{d_n} \int_{x_l}^{\infty} \frac{1}{2\ell_n} \int_{(1-\ell_n)q(u)}^{(1+\ell_n)q(u)} \eta(u, v) dv du = \frac{1}{2\ell_n} (\nu(S_{n,2}^-) - \nu(S_{n,2}^+))$$

with

$$S_{n,2}^{\pm} := \left\{ \frac{d_n}{e_n} (u, (1 \pm \ell_n)v) \mid (u, v) \in S \right\}.$$

Now one can proceed similarly as in (1.9) (using (1.4) and (1.11)) to construct an estimator of $(e_n/d_n) \int_{x_l}^{\infty} q(u)\eta(u, q(u)) du$:

Corollary 2.2. *Let*

$$\hat{S}_{n,2}^{\pm} := \left\{ (u, (1 \pm \ell_n)v) \mid (u, v) \in \frac{n}{ke_n} \hat{T}_n^{\leftarrow}(D_n) \right\}$$

for some sequence $\ell_n \downarrow 0$ such that $k^{-1/2}(w_n(\gamma_1), w_n(\gamma_2)) = o(\ell_n)$. Suppose that all conditions of Theorem 2.1 are fulfilled and, in addition, that an analog to condition (D1) holds where c_n is replaced with $c_n/(1 \pm \ell_n)$. Then

$$\hat{I}_{n,2} := \frac{\hat{\nu}_n(\hat{S}_{n,2}^-) - \hat{\nu}_n(\hat{S}_{n,2}^+)}{2\ell_n} = \frac{e_n}{d_n} \int_{x_l}^{\infty} q(u)\eta(u, q(u)) du (1 + o_P(1)).$$

In a completely analogous way one can estimate $(e_n/d_n) \int_{q(\infty)}^{\infty} q^{\leftarrow}(v)\eta(q^{\leftarrow}(v), v) dv$ by $\hat{I}_{n,1} := (\hat{\nu}_n(\hat{S}_{n,1}^-) - \hat{\nu}_n(\hat{S}_{n,1}^+))/(2\ell_n)$ with

$$\hat{S}_{n,1}^{\pm} := \left\{ ((1 \pm \ell_n)u, v) \mid (u, v) \in \frac{n}{ke_n} \hat{T}_n^{\leftarrow}(D_n) \right\}.$$

Now suppose that both extreme value indices γ_i are positive and that we estimate them by the Hill estimator, i.e., $\hat{\gamma}_1 = k_1^{-1} \sum_{i=1}^{k_1} \log(X_{n-i+1:n}/X_{n-k_1:n})$ with $X_{n-i+1:n}$ denoting the i th largest order statistic among X_1, \dots, X_n , and likewise $\hat{\gamma}_2 = k_2^{-1} \sum_{i=1}^{k_2} \log(Y_{n-i+1:n}/Y_{n-k_2:n})$. It is well known that $k_i^{1/2}(\hat{\gamma}_i - \gamma_i) \rightarrow \mathcal{N}_{(0, \gamma_i^2)}$ if condition (M2) holds and $k_i^{1/2} A_i(n/k_i) \rightarrow 0$. In particular, $\Gamma_i = 0$ if $k = o(k_i)$. However, if $k_i/k \rightarrow \kappa_i \in (0, \infty)$ for both $i = 1$ and $i = 2$, then the joint distribution of Γ_1 and Γ_2 is needed for the construction of confidence intervals.

In the case $k_1 = k_2 = k$, de Haan and Resnick (1993) derived a representation of Γ_i in terms of a Gaussian process under slightly different conditions than used in the present paper. One may mimic their approach to show that under our conditions, $(\Gamma_i/\gamma_i)_{i \in \{1,2\}}$ has the same distribution as $((\int_1^{\infty} t^{-1} W_i(t/\kappa_i) dt - W_i(1/\kappa_i))/\kappa_i)_{i \in \{1,2\}}$ where (W_1, W_2) is a bivariate centered Gaussian process with covariance function given by $Cov(W_1(s), W_1(t)) = \nu((s \vee t, \infty) \times (0, \infty))$, $Cov(W_2(s), W_2(t)) = \nu((0, \infty) \times (s \vee t, \infty))$ and $Cov(W_1(s), W_2(t)) = \nu((s, \infty) \times (t, \infty))$. Direct calculations show that thus $(\Gamma_i/\gamma_i)_{i \in \{1,2\}}$ is a centered Gaussian vector with marginal variances $1/\kappa_i$ and covariance $\nu((\kappa_2, \infty) \times (\kappa_1, \infty))$. Hence, with $z_{1-\alpha/2}$ denoting the standard normal $(1 - \alpha/2)$ -quantile and $\hat{\sigma}^2 := \hat{I}_{n,1}^2/\kappa_1 + \hat{I}_{n,2}^2/\kappa_2 + 2\hat{\nu}_n((\kappa_2, \infty) \times (\kappa_1, \infty))\hat{I}_{n,1}\hat{I}_{n,2}$,

$$\left[\hat{p}_n - k^{-1/2} e_n^{-1} \log c_n \hat{\sigma} z_{1-\alpha/2}, \hat{p}_n + k^{-1/2} e_n^{-1} \log c_n \hat{\sigma} z_{1-\alpha/2} \right] \quad (2.6)$$

is a two-sided confidence interval for p_n with asymptotic confidence level $1 - \alpha$. (This formula is also applicable if one of the κ_i equals ∞ .)

As an alternative to the above approach, one may estimate the density of the spectral measure Φ (cf. Cai et al., 2011) and construct both an estimator for the integrals and for the joint distribution of the limiting random variables on the right-hand side of (2.5) from it.

We conclude this section by indicating how to generalize the main result to \mathbb{R}^d -valued vectors $\mathbf{X}_i = (X_{i,1}, \dots, X_{i,2})$ of arbitrary dimension $d \geq 2$, albeit a detailed discussion is beyond the scope of this paper. An inspection of the proof of Lemma 4.3 reveals that the generalized inverse q^\leftarrow of the function q is used to describe the boundary of the set S as a function of the second coordinate. If $d > 2$ (and hence the generalized inverse is not defined), then an analogous description is needed for all coordinates, i.e. we need d different representations of the set S of the form

$$S = \{\mathbf{x} \in [0, \infty)^d \mid x_i \geq q_i(\mathbf{x}_{-i})\}, \quad 1 \leq i \leq d, \quad (2.7)$$

where $\mathbf{x}_{-i} \in [0, \infty)^{d-1}$ denotes the vector \mathbf{x} with i th coordinate removed and q_i are suitable $[0, \infty]$ -valued functions that are decreasing in each argument. Then one may proceed as in the case $d = 2$ by separately examining the influence of the transformation of each marginal on the ν -measure of the (suitably restricted) set S . Under suitable integrability conditions on the functions q_i and obvious generalizations of the conditions (M1)–(M3), (D1), (D2) and (S1)–(S3), it can be shown that

$$\begin{aligned} & k^{1/2} d_n(\hat{p}_n - p_n) \\ &= \sum_{i=1}^d w_n(\gamma_i) \begin{cases} -\frac{\Gamma_i}{\gamma_i} \int q_i(v) \eta(\tilde{q}_i(v)) 1_{(0, \infty)}(q_i(v)) \mathbb{A}^{d-1}(dv), & \gamma_i > 0 \\ \left(\frac{\alpha_i}{\gamma_i} - \beta - \frac{\Gamma_i}{\gamma_i^2}\right) \int (q_i(v))^{1-\gamma_i} \eta(\tilde{q}_i(v)) 1_{(0, \infty)}(q_i(v)) \mathbb{A}^{d-1}(dv), & \gamma_i < 0 \\ -\Gamma_i \int q_i(v) \eta(\tilde{q}_i(v)) 1_{(0, \infty)}(q_i(v)) \mathbb{A}^{d-1}(dv), & \gamma_i = 0 \end{cases} + o_P(w_n(\gamma_i)) \end{aligned} \quad (2.8)$$

Here \mathbb{A}^{d-1} denotes the Lebesgue measure on $[0, \infty)^{d-1}$ and $\tilde{q}_i(v)$ is the vector in $[0, \infty)^d$ whose i th coordinate equals $q_i(v)$ and the other $d - 1$ coordinates are those of v .

If the boundary ∂S of the set S is sufficiently smooth, then the integrals on the right-hand side of (2.8) can be represented more naturally as integrals w.r.t. certain differential forms (see, e.g., Schreiber, 1977, for an informal introduction to differential forms). More precisely, assume that there exists a set $D \subset [0, \infty)^{d-1}$ and a continuously differentiable function $q : D \rightarrow [0, \infty)$, such that $\partial S = \{(u, q(u)) \mid u \in D\}$ and the map $\Psi : D \rightarrow [0, \infty)^d$, $\Psi(u) = (u, q(u))$ has a Jacobian which is invertible everywhere. Then the right-hand side of (2.8) equals

$$\sum_{i=1}^d w_n(\gamma_i) \begin{cases} -\frac{\Gamma_i}{\gamma_i} \int_{\Psi} pr_i \cdot \eta dx_1 \wedge \dots \wedge dx_{i-1} \wedge dx_{i+1} \wedge \dots \wedge dx_d, & \gamma_i > 0 \\ \left(\frac{\alpha_i}{\gamma_i} - \beta - \frac{\Gamma_i}{\gamma_i^2}\right) \int_{\Psi} (pr_i)^{1-\gamma_i} \cdot \eta dx_1 \wedge \dots \wedge dx_{i-1} \wedge dx_{i+1} \wedge \dots \wedge dx_d, & \gamma_i < 0 \\ -\Gamma_i \int_{\Psi} pr_i \cdot \eta dx_1 \wedge \dots \wedge dx_{i-1} \wedge dx_{i+1} \wedge \dots \wedge dx_d, & \gamma_i = 0 \end{cases} + o_P(w_n(\gamma_i))$$

with pr_i denoting the projection to the i th coordinate. This representation reflects most clearly the fact that the i th term results from the change of the boundary surface of S by the marginal transformation $H_{n,i}$. It is worth mentioning that such a representation can be derived for more general differentiable manifolds ∂S .

3 Analysis of insurance claims

In this section, we discuss issues arising in practical applications. In particular, we will see that, though in the asymptotic setting the constants k and e_n , which are chosen by the statistician, play different roles, from a different perspective essentially only the product ke_n matters.

As an example, we consider a well-known data set of claims to Danish fire insurances. The data set contains losses to building(s), losses to contents and losses to profits (caused by the same fire) observed in the period 01/1980 - 12/2002, discounted to 07/1985. The claims are recorded only if the sum of all components exceeds 1 million Danish Kroner (DKK). Note that due to this recording method, there is an artificial negative dependence between the components, since if one component is smaller than 1 million DKK, the sum of the others must be accordingly larger. To avoid this effect, we therefore consider only those claims for which at least one component exceeds 1 million DKK, which leads to a sample of 3976 claims. Moreover, we focus on the losses to buildings, denoted by X_i as a multiple of one million DKK, and the losses to contents Y_i . A more detailed description of the data can be found in Müller (2008) and Drees and Müller (2008).

As described in the introduction, we assume that a quota reinsurance pays $(1 - \alpha_X)X_i$ for each loss X_i to the building and $(1 - \alpha_Y)Y_i$ for each loss Y_i of content, while an XL-reinsurance pays if the remaining costs $\alpha_X X_i + \alpha_Y Y_i$ exceed a retention level R . We want to estimate the probability $p_n := P(D_n)$ with $D_n := \{\alpha_X X_i + \alpha_Y Y_i > R\}$ that a fire results in a claim to the XL-reinsurance. (More precisely, we estimate the conditional probability given that $\max(X_i, Y_i) > 1$.)

Müller (2008), Section 5.1.2, fitted the following GPD's to the marginal distributions using the Hill estimators based on the $k_1 = 900$ and $k_2 = 600$ largest observations:

$$\hat{F}_i(x) := 1 - \left(1 + \hat{\gamma}_i \frac{x - \mu_i}{\sigma_i}\right)^{-1/\gamma_i}, \quad i = 1, 2, \quad (3.1)$$

with parameters $\hat{\gamma}_1 = 0.57$, $\hat{\sigma}_1 = 0.54$, $\hat{\mu}_1 = 0.91$, $\hat{\gamma}_2 = 0.72$, $\hat{\sigma}_2 = 0.47$ and $\hat{\mu}_2 = 0.15$. (These approximations are sufficient accurate for x satisfying $1 - F_i(x) \leq k_i/n$.) Moreover, he showed that both components of the claim vector are apparently asymptotically dependent.

Note that $\hat{U}_i^{\leftarrow} := 1/(1 - \hat{F}_i)$ can also be interpreted as an estimator $(n/k)T_{n,i}^{\leftarrow}$ for different values of k . However, the number k does not have any operational meaning if one starts with a given approximation of the marginal tails as in (3.1). In that case it seems more natural to reformulate our estimator \hat{p}_n , the main result (2.5) and the resulting confidence interval (2.6) in terms of \hat{U}_i^{\leftarrow} .

First, note that the estimator of the failure probability

$$\hat{p}_n = \frac{1}{e_n} \hat{\nu}_n\left(\frac{n}{ke_n} \hat{T}_n^{\leftarrow}(D_n)\right) = \frac{1}{ke_n} \sum_{i=1}^n \varepsilon_{\hat{U}^{\leftarrow}(X_i, Y_i)}\left(\frac{n}{ke_n} \hat{U}^{\leftarrow}(D_n)\right)$$

depends on the constants k and e_n only via their product ke_n (if the tail estimator \hat{U}^{\leftarrow} is considered fixed). At first glance, this seems peculiar, because in Theorem 2.1 the estimation error seemingly depends on k and e_n in completely different ways. However, note that according to the discussion following Theorem 2.1, for $\gamma_1, \gamma_2 > 0$, approximation (2.5) can be rewritten as

$$\hat{p}_n - p_n = (ke_n)^{-1/2} \log \frac{ke_n}{n} N(1 + o_P(1)) \quad (3.2)$$

for a centered Gaussian random variable N with variance

$$\begin{aligned}\hat{\sigma}_N^2 &= \frac{1}{e_n} \left(\frac{k}{k_1} I_1^2 + \frac{k}{k_2} I_2^2 + 2\nu \left(\left(\frac{k_2}{k}, \infty \right) \times \left(\frac{k_1}{k}, \infty \right) \right) I_1 I_2 \right) \\ &= \frac{k e_n}{k_1} \left(\frac{I_1}{e_n} \right)^2 + \frac{k e_n}{k_2} \left(\frac{I_2}{e_n} \right)^2 + 2\nu \left(\left(\frac{k_2}{k e_n}, \infty \right) \times \left(\frac{k_1}{k e_n}, \infty \right) \right) \frac{I_1}{e_n} \frac{I_2}{e_n},\end{aligned}$$

where $I_1 := (e_n/d_n) \int_{q(\infty)}^\infty q^\leftarrow(v) \eta(q^\leftarrow(v), v) dv$ and $I_2 := (e_n/d_n) \int_{x_l}^\infty q(u) \eta(u, q(u)) du$. Thus I_i/e_n does not depend on e_n , and the distribution of the approximating Gaussian random variable on the right-hand side of (3.2) depends on k and e_n only via their product.

Moreover, also the estimators

$$\begin{aligned}\frac{\hat{I}_{n,1}}{e_n} &= \frac{\hat{\nu}_n(\hat{S}_{n,1}^-) - \hat{\nu}_n(\hat{S}_{n,1}^+)}{2\ell_n e_n} \\ &= \frac{1}{2\ell_n} \cdot \frac{1}{k e_n} \sum_{i=1}^n \varepsilon_{\hat{U}^\leftarrow(X_i, Y_i)} \left(\left\{ ((1 - \ell_n)u, v) \mid (u, v) \in \frac{n}{k e_n} \hat{U}^\leftarrow(D_n) \right\} \setminus \right. \\ &\quad \left. \left\{ ((1 + \ell_n)u, v) \mid (u, v) \in \frac{n}{k e_n} \hat{U}^\leftarrow(D_n) \right\} \right)\end{aligned}$$

and likewise $\hat{I}_{n,2}/e_n$ depend on the product $k e_n$ only. Finally, the covariance term $\nu(k_2/(k e_n), \infty) \times (k_1/(k e_n), \infty) = k^2 e_n / (\lambda k_1 k_2) \nu((k/(\lambda k_1), \infty) \times (k/(\lambda k_2), \infty))$ can be estimated by

$$\frac{k^2 e_n}{\lambda k_1 k_2} \hat{\nu}_n \left(\left(\frac{k}{\lambda k_1}, \infty \right) \times \left(\frac{k}{\lambda k_2}, \infty \right) \right) = \frac{k e_n}{\lambda k_1 k_2} \sum_{i=1}^n \varepsilon_{\hat{U}^\leftarrow(X_i, Y_i)} \left(\left(\frac{n}{\lambda k_1}, \infty \right) \times \left(\frac{n}{\lambda k_2}, \infty \right) \right).$$

Here the choice $\lambda \in (0, 1]$ ensures that \hat{U}^\leftarrow is used only on the range where it is a sufficiently accurate estimator of the true function U^\leftarrow .

To sum up, all estimates only depend on $k e_n$, but not on the numbers k and e_n separately. This product should be chosen as large as possible under the constraints that both marginal approximations of U_i^\leftarrow by \hat{U}_i^\leftarrow and the approximation of the joint distribution of the standardized vector (cf. (1.2)) are reliable. To ensure the former constraint, for the vast majority of the observations (X_i, Y_i) , the indicator of the set $\{\hat{U}^\leftarrow(X_i, Y_i) \in n/(k e_n) \hat{U}^\leftarrow(D_n)\}$ should not depend on the particular values of $\hat{U}_1^\leftarrow(X_i)$ or $\hat{U}_2^\leftarrow(Y_i)$ if these are smaller than n/k_1 or n/k_2 (either because the other component of the vector is so large that the observations lie in the failure set anyway, or because the other component is so small so that the indicator is 0 even if the maximal value n/k_i is attained). To be more concrete, for the failure set $D_n := \{(x, y) \mid \alpha_1 x + \alpha_2 y > R\}$ introduced above, $k e_n$ should be smaller than $\min_{i=1,2} k_i \hat{U}_i^\leftarrow(R/\alpha_i)$, because otherwise for sure $\hat{U}^\leftarrow(x, y) \in (n/k e_n) \hat{U}^\leftarrow(D_n)$ for some values (x, y) for which $\hat{U}^\leftarrow(x, y)$ is not a reliable estimate of $U^\leftarrow(x, y)$.

However, the above crude upper bound $k e_n$ is not sufficient to ensure that \hat{p}_n is a reliable estimate of p_n , because the dependence structure must be accurately described by the exponent measure ν , too. We propose in analogy to the well-known Hill plot, to plot \hat{p}_n versus $k e_n$ and then to choose $k e_n$ in a range where this curve seems stable. In Figure 1 such a graph is shown for the Danish fire insurance data and the failure set $D_n = \{(x, y) \mid x + 0.5y > 100\}$ for values of $k e_n$ ranging from 10^4 to $5 \cdot 10^5$. Note that the aforementioned crude upper bound on $k e_n$ is about $1.7 \cdot 10^6$, but the curve of probability estimates shows a clear downward trend for $k e_n > 2 \cdot 10^5$, which is most

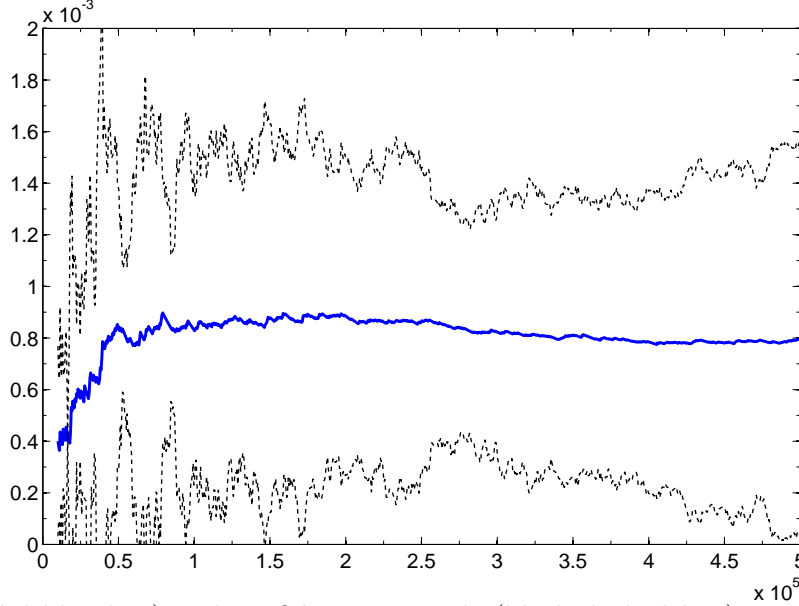


Figure 1: \hat{p}_n (solid blue line) and confidence intervals (black dashed line) versus $k e_n$ for Danish fire insurance claims

likely due to a deviation of the dependence structure from its limit. On the other hand, for values smaller than $5 \cdot 10^4$ the curve is very unstable, too, because the random error is still too large as just a few observations fall into the inflated failure set (e.g., about 25 if $k e_n \approx 3 \cdot 10^4$). This lower bound on $k e_n$ reflects the condition in the asymptotic framework that n is of smaller order than $k_n e_n$ (see condition (S1)). In view of this plot, the choice $k e_n = 2 \cdot 10^5$ seems reasonable.

In addition, Figure 1 shows a two-sided confidence interval with nominal size 0.95 again as a function of $k e_n$. Here we have chosen $\ell_n = 0.1$ and $\lambda = 1$ in the estimator of the variance σ_N^2 described above; other values of λ between $1/2$ and 1 yield essentially the same estimates, while smaller values of ℓ_n lead to larger fluctuations in the confidence bounds, that however are still of a similar size.

For $k e_n = 2 \cdot 10^5$ one obtains a point estimate for p_n of about $8.8 \cdot 10^{-4}$ and a confidence interval $[2.2 \cdot 10^{-4}, 1.54 \cdot 10^{-3}]$. At first glance, this confidence interval seems rather wide. However, one has to be aware of the fact that we estimate the probability of a very rare event which has occurred only twice in the observational period of more than 20 years. Indeed the empirical probability of the event is about $5 \cdot 10^{-4}$, and the Clopper-Pearson confidence interval $[6 \cdot 10^{-5}, 1.8 \cdot 10^{-3}]$ (again with nominal size 0.95) is even wider. It is worth mentioning that both the empirical point estimate and the Clopper-Pearson confidence interval are exactly the same if one wants to estimate the probability that a claim occurs to the XL-reinsurance for any retention level R between 77 and 145 million DKK! Moreover, for retention level above 152 million DKK the point estimate would be 0 and thus useless for purposes of risk management.

4 Proofs

The proof of Theorem 2.1 is based on the decomposition (1.15) of the estimation error. The asymptotic behavior of the leading term $IV + V$ is established in Corollary 4.5. As a preparation

for this result, first we establish an approximation of the random transformation of the marginals defined in (1.12). Thereby we must restrict ourselves to arguments that are neither too small nor too large, which defines a certain subset S_n^* of S . In Lemma 4.2 an upper bound on the difference between the ν -measures of S and S_n^* is derived, while the Lemmas 4.3 and 4.4 analyze the influence of the marginal transformations on the ν -measure of S_n^* .

In Lemma 4.6 an upper bound on the term II of decomposition (1.15) is proved using empirical process theory. Finally, Lemma 4.8 establishes upper bounds on the terms I and III , while Lemma 4.9 takes care of VI .

Lemma 4.1. *Assume that the conditions (M1)–(M3) and (S1) are fulfilled. For $i \in \{1, 2\}$, let $\lambda_{n,i} > 0$ be a decreasing and $\tau_{n,i} < \infty$ an increasing sequence, such that the following conditions are met:*

- (i) $A_i(n/k)(\lambda_{n,i}d_nk/n)^{\rho_i \pm \varepsilon} = o(k^{-1/2}w_n(\gamma_i))$ for some $\varepsilon > 0$
- (ii) If $\gamma_i > 0$, then $k^{-1/2} = o((\lambda_{n,i}d_n/e_n)^{\gamma_i})$.
- (iii) If $\gamma_i < 0$, then $k^{-1/2} = o((\tau_{n,i}d_nk/n)^{\gamma_i})$ and $\log(d_n/e_n) = o((d_nk/n)^{-\gamma_i})$
- (iv) If $\gamma_i = 0$, then $k^{-1/2} \log \tau_{n,i} \rightarrow 0$ and $\log(d_n/e_n) = o(\log c_n)$

Then, for $i \in \{1, 2\}$,

$$\begin{aligned} \frac{d_n}{e_n} H_{n,i}(x) &= T_{n,i}^{\leftarrow} \circ \hat{T}_{n,i} \circ \hat{T}_{n,i}^{(c)\leftarrow} \circ U_i(d_nx) \\ &= \frac{d_n}{e_n} x \left(1 + \begin{cases} -k^{-1/2} \log c_n \left(\frac{\Gamma_i}{\gamma_i} + o_P(1) \right) + O_P(k^{-1/2}(xd_n/e_n)^{-\gamma_i}), & \gamma_i > 0 \\ k^{-1/2}(d_nk/n)^{-\gamma_i} ((\alpha_i/\gamma_i - \beta_i - \Gamma_i/\gamma_i^2 + o_P(1))x^{-\gamma_i} + o_P(1)), & \gamma_i < 0 \\ -k^{-1/2} \log^2 c_n (\Gamma_i/2 + o_P(1)) + O_P(k^{-1/2} \log c_n \log x), & \gamma_i = 0 \end{cases} \right) \end{aligned}$$

uniformly for $x \in [\lambda_{n,i}, \tau_{n,i}]$.

Proof. For notational simplicity, we omit all indices and arguments of the marginal parameters and normalizing functions and their estimators; e.g., we use \hat{a} as a short form of $\hat{a}_i(n/k)$. Moreover, we drop all indices referring to the i th marginal, i.e., we write U instead of U_i , T_n instead of $T_{n,i}$ and so on.

By (2.2), for all $0 < \varepsilon < |\rho|$,

$$\Delta_1(x) := \frac{U(d_nx) - b}{a} - \frac{(xd_nk/n)^\gamma - 1}{\gamma} = O(A(n/k)(xd_nk/n)^{\gamma+\rho \pm \varepsilon}) = o(k^{-1/2}w_n(\gamma)(xd_nk/n)^\gamma) \quad (4.1)$$

uniformly for all $x \geq \lambda_n$, where in the last step we have used condition (i). Now one can conclude that $U(d_nx) \in \hat{T}_n((0, \infty))$ for all $x \in [\lambda_n, \tau_n]$ with probability tending to 1. For example, if $\gamma > 0$, then we have to show that $U(d_nx) > \hat{b} - \hat{a}/\hat{\gamma}$ for all $x \geq \lambda_n$ or, equivalently, (using (M3)) that

$$\Delta_1(\lambda_n) \stackrel{(!)}{>} \frac{\hat{b} - b}{a} - \frac{\hat{a}}{a\hat{\gamma}} - \frac{(\lambda_n d_n k/n)^\gamma - 1}{\gamma} = -\frac{1}{\gamma} \left(\frac{d_n k}{n} \lambda_n \right)^\gamma + O(k^{-1/2})$$

which follows immediately from (4.1), (S1) and (ii).

Hence

$$T_n^{\leftarrow} \circ \hat{T}_n \circ \hat{T}_n^{(c)\leftarrow} \circ U(d_n x) = \left[1 + \frac{\gamma}{a} \left(\hat{a} \frac{(c_n^{-1} (1 + \hat{\gamma} \frac{U(d_n x) - \hat{b}}{\hat{a}})^{1/\hat{\gamma}})^{\hat{\gamma}} - 1}{\hat{\gamma}} + \hat{b} - b \right) \right]^{1/\gamma} =: \tilde{H}(x)$$

if the expression in brackets is strictly positive, which will indeed follow from the calculations below.

Now direct calculations show that

$$\tilde{H}(x) = \left[1 + \gamma \left(c_n^{-\hat{\gamma}} \frac{U(d_n x) - b}{a} + \frac{c_n^{-\hat{\gamma}} - 1}{\hat{\gamma}} \left(\frac{\hat{a}}{a} - \frac{\hat{b} - b}{a} \hat{\gamma} \right) \right) \right]^{1/\gamma}. \quad (4.2)$$

By assumption (M3)

$$\Delta_2 := \frac{\hat{a}}{a} - \frac{\hat{b} - b}{a} \hat{\gamma} - 1 = k^{-1/2} (\alpha - \gamma \beta + o_P(1)). \quad (4.3)$$

If $\gamma > 0$, then the Taylor expansion

$$c_n^{-\hat{\gamma}/\gamma} = c_n^{-1} \left(1 - k^{-1/2} \frac{\Gamma}{\gamma} \log c_n + o_P(k^{-1/2} \log c_n) \right)$$

together with (4.1), (4.3) and (S1) implies

$$\begin{aligned} \tilde{H}(x) &= \left[c_n^{-\hat{\gamma}} \left(\left(\frac{d_n k}{n} x \right)^{\gamma} - 1 + \gamma \Delta_1(x) + \frac{\gamma}{\hat{\gamma}} (1 + \Delta_2) \right) + 1 - \frac{\gamma}{\hat{\gamma}} - \frac{\gamma}{\hat{\gamma}} \Delta_2 \right]^{1/\gamma} \\ &= c_n^{-\hat{\gamma}/\gamma} \frac{d_n k}{n} x \left[1 + O_P(|\Delta_1(x)| + k^{-1/2}) \left(\frac{d_n k}{n} x \right)^{-\gamma} + O_P(k^{-1/2} c_n^{\hat{\gamma}} \left(\frac{d_n k}{n} x \right)^{-\gamma}) \right]^{1/\gamma} \\ &= \frac{d_n}{e_n} x \left(1 - k^{-1/2} \frac{\Gamma}{\gamma} \log c_n + o_P(k^{-1/2} \log c_n) \right) \left[1 + o_P(k^{-1/2} \log c_n) + O_P(k^{-1/2} \left(\frac{d_n}{e_n} x \right)^{-\gamma}) \right], \end{aligned}$$

from which the assertion follows readily.

If $\gamma < 0$, then similar arguments prove

$$\tilde{H}(x) = \frac{d_n}{e_n} x (1 + O_P(k^{-1/2} \log c_n)) \left[1 + k^{-1/2} \frac{1}{\gamma} \left(\alpha - \gamma \beta - \frac{\Gamma}{\gamma} + o_P(1) \right) \left(\frac{d_n k}{n} x \right)^{-\gamma} + o_P(k^{-1/2} w_n(\gamma)) \right],$$

and hence the assertion, because the assumption (iii) ensures that $\log c_n = o(w_n(\gamma))$.

Finally, for $\gamma = 0$, the Taylor expansion

$$c_n^{-\hat{\gamma}} = 1 - \hat{\gamma} \log c_n + \frac{1}{2} \hat{\gamma}^2 \log^2 c_n + O_P(\hat{\gamma}^3 \log^3 c_n)$$

yields

$$\begin{aligned} \tilde{H}(x) &= \exp \left[(1 - \hat{\gamma} \log c_n + O_P(k^{-1} \log^2 c_n)) \left(\log \left(\frac{d_n k}{n} x \right) + \Delta_1(x) \right) + \right. \\ &\quad \left. + \left(-\log c_n + \frac{1}{2} \hat{\gamma} \log^2 c_n + O_P(k^{-1} \log^3 c_n) \right) (1 + \Delta_2) \right] \\ &= \frac{d_n k}{c_n n} x \exp \left[-\hat{\gamma} \log c_n \left(\log c_n + \log \left(\frac{d_n}{e_n} x \right) \right) + O_P(k^{-1/2} \log^2 c_n) + \frac{1}{2} \hat{\gamma} \log^2 c_n \right] \\ &= \frac{d_n}{e_n} x \left[1 - \frac{1}{2} (\Gamma + o_P(1)) k^{-1/2} \log^2 c_n + O_P(k^{-1/2} \log c_n \log x) \right], \end{aligned}$$

which concludes the proof. \square

In what follows we denote by $\lambda_{n,1} \searrow x_l$, $\lambda_{n,2} \searrow q(\infty) := \lim_{x \rightarrow \infty} q(x)$ and $\tau_{n,i} \uparrow \infty$, $i \in \{1, 2\}$, sequences which satisfy the conditions of Lemma 4.1. (These sequences will be specified in the proof of Corollary 4.5.) Note that in particular constant sequences $\lambda_{n,i}, \tau_{n,i} \in (0, \infty)$ satisfy the conditions of Lemma 4.1, provided

$$A_i(n/k)c_n^{\rho_i+\varepsilon} = o(k_n^{-1/2}w_n(\gamma_i)) \quad \text{for } i \in \{1, 2\} \text{ and some } \varepsilon > 0 \quad (4.4)$$

and (S1) holds. Therefore we may and will choose

$$\begin{aligned} \lambda_{n,1} &= x_l & \text{if } x_l := \inf\{x \geq 0 \mid q(x) < \infty\} > 0, \\ \lambda_{n,2} &= q(\infty) & \text{if } q(\infty) > 0, \end{aligned} \quad (4.5)$$

We want to apply the approximations just established to points (x, y) on the boundary of S . To ensure that $x \in [\lambda_{n,1}, \tau_{n,1}]$ and $y \in [\lambda_{n,2}, \tau_{n,2}]$, we consider a subset S_n^* of S that is bounded away from the coordinate axes. More precisely, we define

$$S_n^* := S \cap ([u_n^*, \infty) \times [v_n^*, \infty))$$

with

$$u_n^* := \lambda_{n,1} \vee q^\leftarrow(\tau_{n,2}), \quad v_n^* := \lambda_{n,2} \vee q(\tau_{n,1}).$$

The following lemma implies that the ν -measure of the set $S \setminus S_n^*$ is asymptotically negligible.

Lemma 4.2.

$$\nu(S) - \nu(S_n^*) = O\left(\frac{\lambda_{n,1} - x_l}{q^2(\lambda_{n,1}-)} + \frac{q(\tau_{n,1}) - q(\infty)}{\tau_{n,1}^2} + \frac{\lambda_{n,2} - q(\infty)}{(q^\leftarrow(\lambda_{n,2}))^2} + \frac{q^\leftarrow(\tau_{n,2}) - x_l}{\tau_{n,2}^2}\right)$$

with $q(x-) := \lim_{t \uparrow x} q(t)$.

Proof. First note that $S \subset [x_l, \infty) \times [q(\infty), \infty)$ implies

$$\nu(S) - \nu(S \cap ([0, u_n^*] \times [0, \infty))) \leq \nu([x_l, \lambda_{n,1}] \times [q(\lambda_{n,1}-), \infty)) + \nu([x_l, q^\leftarrow(\tau_{n,2})] \times [\tau_{n,2}, \infty)).$$

The spectral density φ is assumed continuous and hence it is bounded. From (2.4) we conclude that for arbitrary $0 \leq u_0 \leq u_1$ and $v_0 > 0$

$$\nu([u_0, u_1] \times [v_0, \infty)) = O\left(\int_{u_0}^{u_1} \int_{v_0}^{\infty} (u^2 + v^2)^{-3/2} dv du\right) = O\left(\frac{u_1 - u_0}{v_0^2}\right)$$

and thus

$$\nu(S) - \nu(S \cap ([0, u_n^*] \times [0, \infty))) = O\left(\frac{\lambda_{n,1} - x_l}{q^2(\lambda_{n,1}-)} + \frac{q^\leftarrow(\tau_{n,2}) - x_l}{\tau_{n,2}^2}\right).$$

Likewise, one can show that

$$\nu(S \cap ([0, u_n^*] \times [0, \infty))) - \nu(S_n^*) = O\left(\frac{q(\tau_{n,1}) - q(\infty)}{\tau_{n,1}^2} + \frac{\lambda_{n,2} - q(\infty)}{(q^\leftarrow(\lambda_{n,2}))^2}\right).$$

A combination of these two bounds yields the assertion. \square

On the set S_n^* we can now use the approximation from Lemma 4.1 to first examine the influence of the transformation $H_{n,2}$ of the second coordinate on the ν -measure of S_n^* . In a second step we then similarly determine how the ν -measure of this transformed set is altered by the transformation $H_{n,1}$ of the first coordinate. Hereby note that by to Lemma 4.1 the marginal transformations are invertible with probability tending to 1.

Lemma 4.3. *Let $H_n(x, y) := (H_{n,1}(x), H_{n,2}(y)) := \frac{e_n}{d_n} T_n^{\leftarrow} \circ \hat{T}_n \circ \hat{T}_n^{(c)\leftarrow} \circ U(d_n x, d_n y)$. Suppose that the conditions (D2) and (Q1) are met.*

Then one has with $q_n(u) := q(u) \vee v_n^$ and $\tilde{q}_n^{\leftarrow}(v) := q^{\leftarrow}(H_{n,2}^{\leftarrow}(v)) \vee u_n^*$*

$$\begin{aligned}
& \left| \nu(H_n(S_n^*)) - \nu(S_n^*) \right. \\
& \quad \left. + \int_{u_n^*}^{\infty} (H_{n,2}(q_n(u)) - q_n(u)) \eta(u, q_n(u)) du + \int_{H_{n,2}(v_n^*)}^{\infty} (H_{n,1}(\tilde{q}_n^{\leftarrow}(v)) - \tilde{q}_n^{\leftarrow}(v)) \eta(\tilde{q}_n^{\leftarrow}(v), v) dv \right| \\
& = o\left(\int_{u_n^*}^{\infty} |H_{n,2}(q_n(u)) - q_n(u)| \eta(u, q_n(u)) du + \int_{H_{n,2}(v_n^*)}^{\infty} |H_{n,1}(\tilde{q}_n^{\leftarrow}(v)) - \tilde{q}_n^{\leftarrow}(v)| \eta(\tilde{q}_n^{\leftarrow}(v), v) dv \right)
\end{aligned} \tag{4.6}$$

with probability tending to 1.

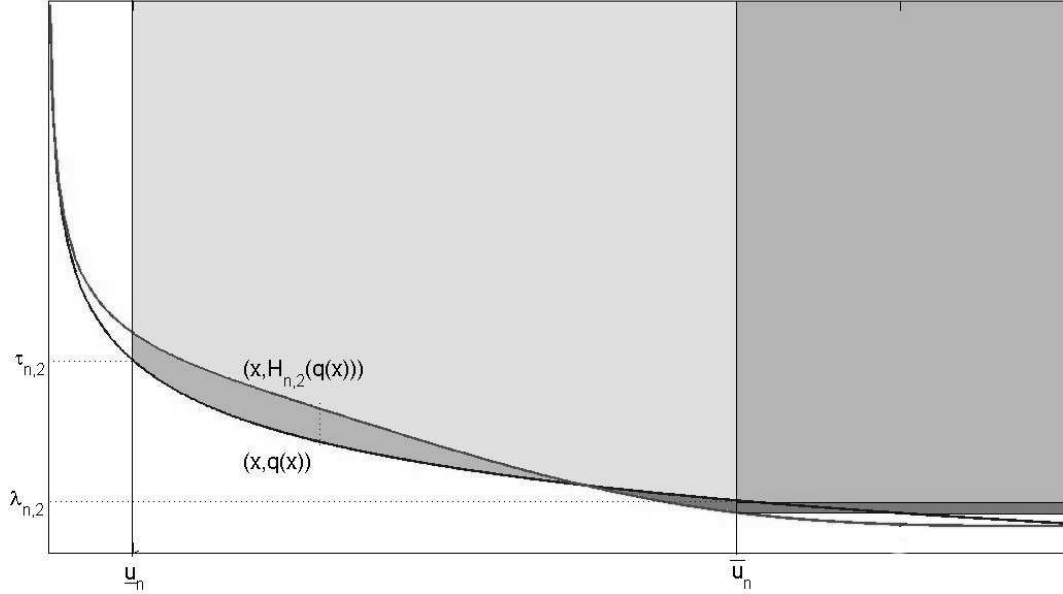


Figure 2: The light and mid grey regions show the approximation S_n^* of the set S , the mid and the dark grey regions the symmetric difference between $\{(x, H_{n,2}(y)) \mid (x, y) \in S_n^*\}$ and S_n^* , where the dark grey region is counted with a positive sign, the mid grey region with a negative sign. (Here it is assumed that $u_n^* = q^{\leftarrow}(\tau_{n,2})$.)

Proof. According to the proof of Lemma 4.1, for all $\delta \in (0, 1)$, on the set $[\lambda_{n,i}(1 - \delta), \tau_{n,i}(1 + \delta)]$ the transformation $H_{n,i}$ is continuous and strictly increasing and $H_{n,i}(x) = x(1 + o(1))$ with probability tending to 1.

We first quantify the influence of the transformation of the second coordinate. Note that

$$\begin{aligned}\nu(S_n^*) &= \int_{u_n^*}^{\infty} \int_{q_n(u)}^{\infty} \eta(u, v) dv du \\ \nu\{(x, H_{n,2}(y)) \mid (x, y) \in S_n^*\} &= \int_{u_n^*}^{\infty} \int_{H_{n,2}(q_n(u))}^{\infty} \eta(u, v) dv du\end{aligned}$$

and hence

$$\nu\{(x, H_{n,2}(y)) \mid (x, y) \in S_n^*\} - \nu(S_n^*) = - \int_{u_n^*}^{\infty} \int_{q_n(u)}^{H_{n,2}(q_n(u))} \eta(u, v) dv du. \quad (4.7)$$

The inner integral equals

$$\begin{aligned}& \int_{q_n(u)}^{H_{n,2}(q_n(u))} \eta(u, v) dv \\ &= (H_{n,2}(q_n(u)) - q_n(u))\eta(u, q_n(u)) + \int_1^{H_{n,2}(q_n(u))/q_n(u)} \left(\frac{\eta(u, q_n(u)w)}{\eta(u, q_n(u))} - 1 \right) dw \eta(u, q_n(u))q_n(u).\end{aligned} \quad (4.8)$$

By the assumptions and Lemma 4.1, $H_{n,2}(q_n(u))/q_n(u) \rightarrow 1$ uniformly for $u \in (u_n^*, \infty)$ as $q_n(u) \in [\lambda_{n,2}, \tau_{n,2}]$ for $u > u_n^*$. Recall that it is assumed that the spectral measure has density φ which is continuous and strictly positive so that φ is uniformly continuous and bounded away from 0. Thus, by (2.4),

$$\frac{\eta(u, vw)}{\eta(u, v)} = \left(\frac{u^2 + v^2}{u^2 + v^2 w^2} \right)^{3/2} \frac{\varphi(\arctan \frac{vw}{u})}{\varphi(\arctan \frac{v}{u})} \rightarrow 1$$

as $w \rightarrow 1$ uniformly for $u, v > 0$, since

$$\frac{u^2 + v^2}{u^2 + v^2 w^2} = 1 + \frac{1 - w^2}{(u/v)^2 + w^2} \rightarrow 1$$

and

$$\left| \frac{\varphi(\arctan \frac{vw}{u})}{\varphi(\arctan \frac{v}{u})} - 1 \right| \leq \frac{|\varphi(\arctan \frac{vw}{u}) - \varphi(\arctan \frac{v}{u})|}{\inf_{0 \leq \vartheta \leq \pi/2} \varphi(\vartheta)} \rightarrow 0$$

as $w \rightarrow 1$ uniformly for $u, v > 0$ by the uniform continuity of φ and

$$\left| \arctan \frac{vw}{u} - \arctan \frac{v}{u} \right| = \frac{v/u}{1 + ((v/u)(1 + \theta(w-1)))^2} |w-1| \leq \sup_{z>0} \frac{z}{1 + z^2/4} |w-1|$$

for $w > 1/2$ and some $\theta \in (0, 1)$, which holds by the mean value theorem.

Therefore,

$$\int_1^{H_{n,2}(q_n(u))/q_n(u)} \left(\frac{\eta(u, q_n(u)w)}{\eta(u, q_n(u))} - 1 \right) dw = o(H_{n,2}(q_n(u))/q_n(u) - 1)$$

which, combined with (4.7) and (4.8), yields

$$\begin{aligned}& \nu\{(x, H_{n,2}(y)) \mid (x, y) \in S_n^*\} - \nu(S_n^*) + \int_{u_n^*}^{\infty} (H_{n,2}(q_n(u)) - q_n(u))\eta(u, q_n(u)) du \\ &= o\left(\int_{u_n^*}^{\infty} |H_{n,2}(q_n(u)) - q_n(u)| \eta(u, q_n(u)) du \right).\end{aligned} \quad (4.9)$$

One can derive an analogous approximation of the difference between $\nu\{(x, H_{n,2}(y)) \mid (x, y) \in S\}$ and $\nu\{(H_{n,1}(x), H_{n,2}(y)) \mid (x, y) \in S\}$ by similar arguments if one interchanges the order of integration:

$$\begin{aligned} & \left| \nu(H_n(S_n^*)) - \nu\{(x, H_{n,2}(y)) \mid (x, y) \in S_n^*\} + \int_{H_{n,2}(v_n^*)}^{\infty} (H_{n,1}(\tilde{q}_n^{\leftarrow}(v)) - \tilde{q}_n^{\leftarrow}(v)) \eta(\tilde{q}_n^{\leftarrow}(v), v) dv \right| \\ &= o\left(\int_{H_{n,2}(v_n^*)}^{\infty} |H_{n,1}(\tilde{q}_n^{\leftarrow}(v)) - \tilde{q}_n^{\leftarrow}(v)| \eta(\tilde{q}_n^{\leftarrow}(v), v) dv \right). \end{aligned} \quad (4.10)$$

Summing up (4.9) and (4.10), we arrive at the assertion. \square

In the next lemma, we calculate the limits of the integrals arising in Lemma 4.3 using the approximation established in Lemma 4.1.

Lemma 4.4. *Suppose that the conditions of Lemma 4.3 and, in addition, the following conditions are fulfilled for some $x_0 \in (x_l, q^{\leftarrow}(q(\infty)))$, $y_0 \in (q(\infty), q(x_l))$:*

$$\int_{y_0}^{\infty} (q^{\leftarrow}(v))^{1-\gamma_1} v^{-3} dv < \infty \quad \text{or} \quad \lambda_{n,1}^{1-\gamma_1} = o(\log c_n) \quad (4.11)$$

$$\int_{x_0}^{\infty} (q(u))^{1-\gamma_2} u^{-3} du < \infty \quad \text{or} \quad \lambda_{n,2}^{1-\gamma_2} = o(\log c_n) \quad (4.12)$$

Then the following approximations hold true:

(i)

$$\begin{aligned} & \frac{k^{1/2}}{w_n(\gamma_2)} \int_{u_n^*}^{\infty} (H_{n,2}(q_n(u)) - q_n(u)) \eta(u, q_n(u)) du \\ & \rightarrow \begin{cases} -\frac{\Gamma_2}{\gamma_2} \int_{x_l}^{\infty} q(u) \eta(u, q(u)) du, & \gamma_2 > 0 \\ \left(\frac{\alpha_2}{\gamma_2} - \beta_2 - \frac{\Gamma_2}{\gamma_2^2}\right) \int_{x_l}^{\infty} (q(u))^{1-\gamma_2} \eta(u, q(u)) du, & \gamma_2 < 0 \\ -\Gamma_2 \int_{x_l}^{\infty} q(u) \eta(u, q(u)) du, & \gamma_2 = 0 \end{cases} \end{aligned}$$

Moreover,

$$\int_{u_n^*}^{\infty} |H_{n,2}(q_n(u)) - q_n(u)| \eta(u, q_n(u)) du = O(k^{-1/2} w_n(\gamma_2)).$$

(ii)

$$\begin{aligned} & \frac{k^{1/2}}{w_n(\gamma_1)} \int_{H_{n,2}(v_n^*)}^{\infty} (H_{n,1}(\tilde{q}_n^{\leftarrow}(v)) - \tilde{q}_n^{\leftarrow}(v)) \eta(\tilde{q}_n^{\leftarrow}(v), v) dv \\ & \rightarrow \begin{cases} -\frac{\Gamma_1}{\gamma_1} \int_{q(\infty)}^{\infty} q^{\leftarrow}(v) \eta(q^{\leftarrow}(v), v) dv, & \gamma_1 > 0 \\ \left(\frac{\alpha_1}{\gamma_1} - \beta_1 - \frac{\Gamma_1}{\gamma_1^2}\right) \int_{q(\infty)}^{\infty} (q^{\leftarrow}(v))^{1-\gamma_1} \eta(q^{\leftarrow}(v), v) dv, & \gamma_1 < 0 \\ -\Gamma_1 \int_{q(\infty)}^{\infty} q^{\leftarrow}(v) \eta(q^{\leftarrow}(v), v) dv, & \gamma_1 = 0 \end{cases} \end{aligned}$$

Furthermore,

$$\int_{H_{n,2}(v_n^*)}^{\infty} |H_{n,1}(\tilde{q}_n^{\leftarrow}(v)) - \tilde{q}_n^{\leftarrow}(v)| \eta(\tilde{q}_n^{\leftarrow}(v), v) dv = O(k^{-1/2} w_n(\gamma_1)).$$

Proof. ad (i): Because the spectral density φ is bounded, there exists a constant $K > 0$ such that

$$\eta(u, q(u)) \leq K(u^2 + (q(u))^2)^{-3/2} \leq K(u^{-3} \wedge (q(u))^{-3}) \quad \forall u > 0. \quad (4.13)$$

Hence $q_n(u)\eta(u, q_n(u)) \leq K(q(u))^{-2}$ for $u \in [x_l, x_0]$ and n sufficiently large and $q_n(u)\eta(u, q_n(u)) \leq Kq(x_0)u^{-3}$ for $u > x_0$. Therefore,

$$\lim_{n \rightarrow \infty} \int_{u_n^*}^{\infty} q_n(u)\eta(u, q_n(u)) du = \int_{x_l}^{\infty} q(u)\eta(u, q(u)) du < \infty \quad (4.14)$$

by the dominated convergence theorem and $u_n^* \downarrow x_l$.

Now, we distinguish three cases.

If $\gamma_2 > 0$, then by Lemma 4.1 and $d_n \asymp e_n$

$$\begin{aligned} \frac{k^{1/2}}{\log c_n} \int_{u_n^*}^{\infty} (H_{n,2}(q_n(u)) - q_n(u))\eta(u, q_n(u)) du &= -\left(\frac{\Gamma_2}{\gamma_2} + o_P(1)\right) \int_{u_n^*}^{\infty} q_n(u)\eta(u, q_n(u)) du \\ &\quad + O_P\left(\frac{1}{\log c_n} \int_{u_n^*}^{\infty} (q_n(u))^{1-\gamma_2}\eta(u, q_n(u)) du\right). \end{aligned}$$

Because of (4.13) and (4.12)

$$\int_{u_n^*}^{\infty} (q_n(u))^{1-\gamma_2}\eta(u, q_n(u)) du \leq K(q(x_0))^{-2-\gamma_2}(x_0-x_l) + K \int_{x_0}^{\infty} (q(u) \vee \lambda_{n,2})^{1-\gamma_2} u^{-3} du = o(\log c_n). \quad (4.15)$$

Hence, in view of (4.14), we have

$$\int_{u_n^*}^{\infty} (H_{n,2}(q_n(u)) - q_n(u))\eta(u, q_n(u)) du = -k^{-1/2} \log c_n \frac{\Gamma_2}{\gamma_2} \int_{x_l}^{\infty} q(u)\eta(u, q(u)) du + o_P(k^{-1/2} \log c_n).$$

If $\gamma_2 < 0$, then the assertion follows similarly from Lemma 4.1 and (4.13).

Finally, in the case $\gamma_2 = 0$

$$\int_{x_l}^{\infty} q_n(u) |\log q_n(u)| \eta(u, q_n(u)) du \leq K \sup_{x \leq x_0} \frac{|\log q(x)|}{(q(x))^2} + K \sup_{x \geq x_0} q(x) |\log q(x)| \int_{x_0}^{\infty} u^{-3} du < \infty. \quad (4.16)$$

Hence, similarly as in the first case, we may conclude the assertion from Lemma 4.1.

ad (ii): The second assertion can be proved in a very similar fashion using $q(H_{n,2}(v_n^*)) \rightarrow q(\infty)$ and the fact that $\tilde{q}_n^{\leftarrow}(u) \rightarrow q^{\leftarrow}(u)$ for Lebesgue-almost all $u > q(\infty)$, because of Lemma 4.1 and the Lebesgue-almost surely continuity of q^{\leftarrow} . For that reason, we only give the analog to the bound (4.15) for the integral under consideration in the case $\gamma_1 > 0$.

For $y_0 \in (q(\infty), q(x_l))$ and all sufficiently large n , we have

$$\int_{H_{n,2}(v_n^*)}^{y_0} (\tilde{q}_n^{\leftarrow}(v))^{1-\gamma_1} \eta(\tilde{q}_n^{\leftarrow}(v), v) dv \leq K(\tilde{q}_n^{\leftarrow}(y_0))^{-2-\gamma_1} (y_0 - q(\infty)) = O(1).$$

If $\gamma_1 \leq 1$, then

$$\int_{y_0}^{\infty} (\tilde{q}_n^{\leftarrow}(v))^{1-\gamma_1} \eta(\tilde{q}_n^{\leftarrow}(v), v) dv \leq K(\tilde{q}_n^{\leftarrow}(y_0))^{1-\gamma_1} \int_{y_0}^{\infty} v^{-3} dv = O(1).$$

Finally, if $\gamma_1 > 1$, then by the monotonicity of q^\leftarrow and the asymptotic behavior of $H_{n,2}$ we have for all $\delta > 0$ and sufficiently large n

$$\begin{aligned} \int_{y_0}^{\infty} (\tilde{q}_n^\leftarrow(v))^{1-\gamma_1} \eta(\tilde{q}_n^\leftarrow(v), v) dv &\leq K \int_{y_0}^{\infty} \left((q^\leftarrow(v(1+\delta)))^{1-\gamma_1} \wedge \lambda_{n,1}^{1-\gamma_1} \right) v^{-3} dv \\ &= O\left(\int_{y_0(1+\delta)}^{\infty} (q^\leftarrow(v))^{1-\gamma_1} v^{-3} dv \wedge \lambda_{n,1}^{1-\gamma_1} \right) = o(\log c_n) \end{aligned}$$

by condition (4.11). \square

The following result gives sufficient conditions such that the difference between the ν -measure of S and of the truncated set after the marginal transformations (i.e. $d_n(IV + V)$ in (1.15)) can be approximated by the limiting terms in Lemma 4.4. For the sake of simplicity, we assume that d_n and e_n are of the same order, but it is not difficult to prove similar results under weaker conditions on d_n/e_n . Moreover, one can weaken the condition (S2) and the assumptions (Q2) could be replaced with rather strong conditions on the rate at which k tends to ∞ .

Corollary 4.5. *If the conditions (M1)–(M3), (D2), (Q1), (Q2) and (S1)–(S3) are fulfilled, then*

$$\begin{aligned} &\nu(H_n(S_n^*)) - \nu(S) \\ &= k^{-1/2} w_n(\gamma_1) \begin{cases} -\frac{\Gamma_1}{\gamma_1} \int_{q(\infty)}^{\infty} q^\leftarrow(v) \eta(q^\leftarrow(v), v) dv, & \gamma_1 > 0 \\ \left(\frac{\alpha_1}{\gamma_1} - \beta_1 - \frac{\Gamma_1}{\gamma_1^2}\right) \int_{q(\infty)}^{\infty} (q^\leftarrow(v))^{1-\gamma_1} \eta(q^\leftarrow(v), v) dv, & \gamma_1 < 0 \\ -\Gamma_1 \int_{q(\infty)}^{\infty} q^\leftarrow(v) \eta(q^\leftarrow(v), v) dv, & \gamma_1 = 0 \end{cases} \\ &\quad + k^{-1/2} w_n(\gamma_2) \begin{cases} -\frac{\Gamma_2}{\gamma_2} \int_{x_l}^{\infty} q(u) \eta(u, q(u)) du, & \gamma_2 > 0 \\ \left(\frac{\alpha_2}{\gamma_2} - \beta_2 - \frac{\Gamma_2}{\gamma_2^2}\right) \int_{x_l}^{\infty} (q(u))^{1-\gamma_2} \eta(u, q(u)) du, & \gamma_2 < 0 \\ -\Gamma_2 \int_{x_l}^{\infty} q(u) \eta(u, q(u)) du, & \gamma_2 = 0 \end{cases} \\ &\quad + o_P(k^{-1/2}(w_n(\gamma_1) \vee w_n(\gamma_2))) \end{aligned} \tag{4.17}$$

Proof. In view of the Lemmas 4.2 – 4.4, it suffices to define sequences $\lambda_{n,i}$ and $\tau_{n,i}$, $i \in \{1, 2\}$, such that the conditions (i)–(iv) of Lemma 4.1 and (4.11) and (4.12) are fulfilled and

$$\frac{\lambda_{n,1} - x_l}{q^2(\lambda_{n,1}-)} + \frac{q(\tau_{n,1}) - q(\infty)}{\tau_{n,1}^2} = o(k^{-1/2} w_n(\gamma_1)), \quad \frac{\lambda_{n,2} - q(\infty)}{(q^\leftarrow(\lambda_{n,2}))^2} + \frac{q^\leftarrow(\tau_{n,2}) - x_l}{\tau_{n,2}^2} = o(k^{-1/2} w_n(\gamma_2)).$$

Note that we can check these conditions for $i = 1$ and $i = 2$ separately. We focus on the sequences $\lambda_{n,1}$ and $\tau_{n,1}$, since the case $i = 2$ can be treated analogously if x_l is replaced with $q(\infty)$ and q with q^\leftarrow . Again we distinguish three cases depending on the sign of γ_1 .

If $\gamma_1 > 0$, then $\tau_{n,1}$ must only satisfy $(q(\tau_{n,1}) - q(\infty))/\tau_{n,1}^2 = o(k^{-1/2} \log c_n)$, which can easily be fulfilled by letting $\tau_{n,1}$ tend to ∞ sufficiently fast.

The sequence $\lambda_{n,1}$ has to satisfy the conditions (i) and (ii) of Lemma 4.1, (4.11) and $(\lambda_{n,1} - x_l)/q^2(\lambda_{n,1}) = o(k^{-1/2} \log c_n)$. If $x_l > 0$, then $\lambda_{n,1} = x_l$ does the job, because condition (i) of Lemma 4.1 is implied by (S2).

If $x_l = 0$ and $\gamma_1 \leq 1$, then the integrability condition of (4.11) is trivial. Moreover, $\lambda_{n,1} := k^{-1/2}(\log c_n)^{1/2} \rightarrow 0$ obviously fulfills 4.1 (ii) and $(\lambda_{n,1} - x_l)/q^2(\lambda_{n,1}) = O(\lambda_{n,1}) = o(k^{-1/2} \log c_n)$. Condition 4.1 (i) follows from (S2) and (S3), which implies $c_n \lambda_{n,1} \rightarrow \infty$.

Finally, if $x_l = 0$ and $\gamma_1 > 1$, then $\lambda_{n,1} := (k^{-1/2} \log c_n)^{1/\gamma_1}$ fulfills 4.1 (ii), 4.1 (i) follows from (S2) and (S3) as above, and (Q2) implies

$$\frac{\lambda_{n,1} - x_l}{q^2(\lambda_{n,1})} = O\left(\frac{\lambda_{n,1}^{\gamma_1}}{|\log \lambda_{n,1}|^2}\right) = O\left(k^{-1/2} \frac{\log c_n}{|\log(k^{-1/2} \log c_n)|^2}\right) = o(k^{-1/2} \log c_n)$$

by (S1). Furthermore, the integrability condition of (4.11) is fulfilled, because (Q2) implies $(v/\log v)^{2/(1-\gamma_1)} = O(q^{\leftarrow}(v))$ as $v \rightarrow \infty$.

Next we consider the case $-1/2 < \gamma_1 < 0$, when the integrability condition of (4.11) is trivial. If $x_l > 0$, then we can argue as above that $\lambda_{n,1} = x_l$ satisfies all conditions on $\lambda_{n,1}$. If $x_l = 0$, then define $\lambda_{n,1} = c_n^{-1} \varphi_n$ for some $\varphi_n \rightarrow \infty$ sufficiently slowly, so that 4.1 (i) follows from (S2). Further $(\lambda_{n,1} - x_l)/q^2(\lambda_{n,1}) = O(c_n^{-1} \varphi_n) = o(k^{-1/2} c_n^{-\gamma_1})$ follows from assumption (S3).

The conditions on $\tau_{n,1}$ read as $(q(\tau_{n,1}) - q(\infty))/\tau_{n,1}^2 = o(k^{-1/2} c_n^{-\gamma_1})$ and $k^{-1/2} = o((c_n \tau_{n,1})^{\gamma_1})$ in this case, which are fulfilled by $\tau_{n,1} = k^{1/2} c_n^{\gamma_1} \rightarrow \infty$.

In the case $\gamma_1 = 0$ the integrability condition of (4.11) is again trivial and $\lambda_{n,1} = x_l$ if $x_l > 0$, and $\lambda_{n,1} = c_n^{-1} \log c_n$ if $x_l = 0$ does the job. Moreover, it is easily checked that $\tau_{n,1} = k^{1/4}$ satisfies $(q(\tau_{n,1}) - q(\infty))/\tau_{n,1}^2 = o(k^{-1/2} \log^2 c_n)$ and condition 4.1 (iv). \square

Observe that we have verified stronger conditions on $\lambda_{n,1}$ and $\tau_{n,1}$ than actually necessary, if $w_n(\gamma_1) = o(w_n(\gamma_2))$. A refined analysis would lead to weaker, but more complex conditions on q and k that depend on both the values of γ_1 and γ_2 at the same time. (Also the proof would become more lengthy as one had to consider 9 cases arising from different combinations of signs of γ_1 and γ_2 .) Moreover, note that for the above choice of $\lambda_{n,i}$ one has

$$c_n \lambda_{n,i} \rightarrow \infty, \quad i \in \{1, 2\}, \quad (4.18)$$

and

$$\lambda_{n,i}^{-\gamma_i} = O(k^{1/2} / \log c_n) \quad \text{if } \gamma_i > 0, \quad i \in \{1, 2\}. \quad (4.19)$$

Now we use classical empirical process theory to establish a uniform bound on $\nu_n(B) - E\nu_n(B)$ and thus on term II in decomposition (1.15).

Lemma 4.6. *Under the conditions of Theorem 2.1, one has*

$$\nu_n(B) - E\nu_n(B)|_{B=(d_n/e_n)H_n(S_n^*)} = o_P(k^{-1/2}(w_n(\gamma_1) \vee w_n(\gamma_2))).$$

Proof. Note that by (4.2) one has

$$\frac{d_n}{e_n} H_n(x_1, x_2) = (\tilde{H}_{\vartheta_1, \chi_1, \xi_1}^{(n,1)}(x_1), \tilde{H}_{\vartheta_2, \chi_2, \xi_2}^{(n,2)}(x_2))$$

for $(x_1, x_2) \in [u_n^*, \infty) \times [v_n^*, \infty)$ with $\tilde{H}_{\vartheta_i, \chi_i, \xi_i}^{(n,i)}$ defined by (2.3) and

$$-\vartheta_i = \chi_i = k^{1/2}(\hat{\gamma}_i - \gamma_i), \quad \xi_i = k^{1/2}\left(\frac{\hat{a}_i(n/k)}{a_i(n/k)} - 1 - \frac{\hat{b}_i(n/k) - b_i(n/k)}{a_i(n/k)} \hat{\gamma}_i\right).$$

Since, according to condition (M3), these random variables are stochastically bounded, it suffices to prove that for all $M > 0$

$$\sup_{\max(|\vartheta_i|, |\chi_i|, |\xi_i|) \leq M} \left| \nu_n(E_{(\vartheta_i, \chi_i, \xi_i)_{i=1,2}}^{(n)}) - E\nu_n(E_{(\vartheta_i, \chi_i, \xi_i)_{i=1,2}}^{(n)}) \right| = o_P(k^{-1/2}(w_n(\gamma_1) \vee w_n(\gamma_2)))$$

where

$$E_{(\vartheta_i, \chi_i, \xi_i)_{i=1,2}}^{(n)} := \{(\tilde{H}_{\vartheta_1, \chi_1, \xi_1}^{(n,1)}(x_1), \tilde{H}_{\vartheta_2, \chi_2, \xi_2}^{(n,2)}(x_2)) \mid (x_1, x_2) \in S_n^*\}.$$

Letting $\theta := (\vartheta_i, \chi_i, \xi_i)_{i=1,2}$ and

$$Z_n(\theta) := \frac{k^{1/2}}{w_n(\gamma_1) \vee w_n(\gamma_2)} (\nu_n(E_\theta^{(n)}) - E\nu_n(E_\theta^{(n)})), \quad \theta \in [-M, M]^6,$$

we have to prove that Z_n tends to 0 in probability uniformly. To this end, we establish asymptotic equicontinuity of Z_n , i.e.

$$\lim_{\delta \downarrow 0} \limsup_{n \rightarrow \infty} P \left\{ \sup_{\theta, \psi \in [-M, M]^6, \|\theta - \psi\|_\infty \leq \delta} |Z_n(\theta) - Z_n(\psi)| > \eta \right\} = 0 \quad \forall \eta > 0, \quad (4.20)$$

and convergence in probability of $Z_n(\theta)$ for all $\theta \in [-M, M]^6$ (see van der Vaart and Wellner, 2000, Theorem 1.5.7).

For the proof of asymptotic equicontinuity, it is crucial that the functions $\tilde{H}_{\vartheta_i, \chi_i, \xi_i}^{(n,i)}(x_i)$ are decreasing in all three parameters for all $(x_1, x_2) \in [u_n^*, \infty) \times [v_n^*, \infty)$. For ξ_i resp. ϑ_i this monotonicity is an immediate consequence of the facts that $(c_n^{-\gamma} - 1)/\gamma$ is negative and increasing in γ (for $c_n > 1$) and that $(1 + \gamma t)^{1/\gamma}$ is increasing in t . Because $c_n^{-\gamma}$ is a decreasing function of γ , the monotonicity in χ_i follows from (2.2), (4.18) and condition (i) of Lemma 4.1, which imply

$$\begin{aligned} \frac{U_i(d_n x) - b_i(n/k)}{a_i(n/k)} &= \frac{(x_i d_n k/n)^{\gamma_i} - 1}{\gamma_i} + O(A_i(n/k)(x_i d_n k/n)^{\gamma_i + \rho_i + \varepsilon}) \\ &= \frac{(x_i c_n d_n / e_n)^{\gamma_i} - 1}{\gamma_i} + o((c_n x_i)^{\gamma_i} k^{-1/2} w_n(\gamma_i)) \\ &> 0 \end{aligned}$$

for sufficiently large n .

The monotonicity of $H_{\cdot, \cdot, \cdot}^{(n,i)}(x_i)$ implies that the sets $E_{(\vartheta_i, \chi_i, \xi_i)_{i=1,2}}^{(n)}$ are increasing in all parameters. Hence, for arbitrary $\theta, \psi \in [-M, M]^6$

$$|Z_n(\theta) - Z_n(\psi)| \leq \frac{k^{1/2}}{w_n(\gamma_1) \vee w_n(\gamma_2)} (\nu_n(E_{\theta \vee \psi}^{(n)} \setminus E_{\theta \wedge \psi}^{(n)}) + E\nu_n(E_{\theta \vee \psi}^{(n)} \setminus E_{\theta \wedge \psi}^{(n)}))$$

where $\theta \vee \psi$ resp. $\theta \wedge \psi$ denote the coordinatewise maximum resp. minimum of θ and ψ .

To establish asymptotic equicontinuity of Z_n , we cover the parameter space $[-M, M]^6$ with hypercubes $I_l := \times_{i=1}^6 [l_i \delta, (l_i + 1)\delta]$, $-\lceil M/\delta \rceil \leq l_i \leq \lfloor M/\delta \rfloor$, for some small $\delta > 0$ (depending on the value η in (4.20)) to be specified later on. For $\theta, \psi \in [-M, M]^6$ with $\|\theta - \psi\|_\infty \leq \delta$ and $l(\theta) := (\lfloor \theta_i/\delta \rfloor)_{1 \leq i \leq 6}$, one has $\|l(\theta) - l(\psi)\| \leq 1$ and thus

$$\begin{aligned} &|Z_n(\theta) - Z_n(\psi)| \\ &\leq |Z_n(\theta) - Z_n(l(\theta)\delta)| + |Z_n(\psi) - Z_n(l(\psi)\delta)| + |Z_n(l(\theta)\delta) - Z_n(l(\psi)\delta)| \\ &\leq 3 \max_{l \in \{-\lceil M/\delta \rceil, \dots, \lfloor M/\delta \rfloor\}^6} \sup_{t, u \in I_l} |Z_n(t) - Z_n(u)| \\ &\leq 3 \frac{k^{1/2}}{w_n(\gamma_1) \vee w_n(\gamma_2)} \max_{l \in \{-\lceil M/\delta \rceil, \dots, \lfloor M/\delta \rfloor\}^6} (\nu_n(E_{(l+1)\delta}^{(n)} \setminus E_{l\delta}^{(n)}) + E\nu_n(E_{(l+1)\delta}^{(n)} \setminus E_{l\delta}^{(n)})) \quad (4.21) \end{aligned}$$

where $(l+1)\delta := ((l_i+1)\delta)_{1 \leq i \leq 6}$. By (D1), the expectation can be approximated as follows:

$$E\nu_n(E_{(l+1)\delta}^{(n)} \setminus E_{l\delta}^{(n)}) = \frac{n}{k} P\{T_n^{\leftarrow}(X, Y) \in E_{(l+1)\delta}^{(n)} \setminus E_{l\delta}^{(n)}\} = \nu(E_{(l+1)\delta}^{(n)} \setminus E_{l\delta}^{(n)}) + O(A_0(n/k)). \quad (4.22)$$

To bound the right-hand side, first note that by similar calculations as in the proof of Lemma 4.1, one obtains

$$\begin{aligned} & \tilde{H}_{\vartheta_i, \chi_i, \xi_i}^{(n,i)}(x) \\ &= \frac{d_n}{e_n} x \left(1 + \begin{cases} -k^{-1/2} \log c_n(\frac{\chi_i}{\gamma_i} + o_P(1)) + O_P(k^{-1/2}(x d_n/e_n)^{-\gamma_i}), & \gamma_i > 0 \\ k^{-1/2}(d_n k/n)^{-\gamma_i} ((\xi_i/\gamma_i + \vartheta_i/\gamma_i^2 + o_P(1))x^{-\gamma_i} + o_P(1)), & \gamma_i < 0 \\ -k^{-1/2} \log^2 c_n(\chi_i + \vartheta_i/2 + o_P(1)) + O_P(k^{-1/2} \log c_n \log x), & \gamma_i = 0 \end{cases} \right) \end{aligned}$$

uniformly for $x \in [\lambda_{n,i}, \tau_{n,i}]$. That means that under the same conditions as in Lemma 4.1 one can prove an analogous approximation where Γ_i is replaced with χ_i if $\gamma_i > 0$, $\alpha_i/\gamma_i - \beta_i - \Gamma_i/\gamma_i^2$ is replaced with $\xi_i/\gamma_i + \vartheta_i/\gamma_i^2$ if $\gamma_i < 0$, and Γ_i is replaced with $2\chi_i + \vartheta_i$ in the case $\gamma_i = 0$. Hence, we may also conclude a corresponding analog to Corollary 4.5, i.e. $\nu((e_n/d_n)E_{(\vartheta_i, \chi_i, \xi_i)_{i=1,2}}^{(n)}) - \nu(S)$ equals the right-hand side of (4.17) with the above substitutions. Because all integrals are finite, there exists a constant $K > 0$ such that for sufficiently large n

$$\nu(E_{(l+1)\delta}^{(n)}) - \nu(E_{l\delta}^{(n)}) \leq \frac{e_n}{d_n} K \delta k^{-1/2} (w_n(\gamma_1) \vee w_n(\gamma_2))$$

uniformly for all $l \in \{-\lceil M/\delta \rceil \dots \lfloor M/\delta \rfloor\}^6$. A combination with (4.22), $e_n \asymp d_n$ and condition (S2) shows that to each $\eta > 0$ there exists $\delta > 0$ such that for sufficiently large n

$$E\nu_n(E_{(l+1)\delta}^{(n)} \setminus E_{l\delta}^{(n)}) \leq \frac{\eta}{12} k^{-1/2} (w_n(\gamma_1) \vee w_n(\gamma_2)). \quad (4.23)$$

In view of (4.21), we obtain

$$\begin{aligned} & P\left\{ \sup_{\theta, \psi \in [-M, M]^6, \|\theta - \psi\|_\infty \leq \delta} |Z_n(\theta) - Z_n(\psi)| > \eta \right\} \\ & \leq P\left\{ \max_{l \in \{-\lceil M/\delta \rceil \dots \lfloor M/\delta \rfloor\}^6} (\nu_n(E_{(l+1)\delta}^{(n)} \setminus E_{l\delta}^{(n)}) + E\nu_n(E_{(l+1)\delta}^{(n)} \setminus E_{l\delta}^{(n)})) \right. \\ & \quad \left. > \frac{\eta}{3} k^{-1/2} (w_n(\gamma_1) \vee w_n(\gamma_2)) \right\} \\ & \leq \sum_{l \in \{-\lceil M/\delta \rceil \dots \lfloor M/\delta \rfloor\}^6} P\left\{ |\nu_n(E_{(l+1)\delta}^{(n)} \setminus E_{l\delta}^{(n)}) - E\nu_n(E_{(l+1)\delta}^{(n)} \setminus E_{l\delta}^{(n)})| \right. \\ & \quad \left. > \frac{\eta}{6} k^{-1/2} (w_n(\gamma_1) \vee w_n(\gamma_2)) \right\}. \end{aligned}$$

Therefore the asserted asymptotic equicontinuity (4.20) follows from (4.23) and Chebyshev's inequality applied to the binomial random variables $k\nu_n(E_{(l+1)\delta}^{(n)} \setminus E_{l\delta}^{(n)})$:

$$\begin{aligned} & P\left\{ |\nu_n(E_{(l+1)\delta}^{(n)} \setminus E_{l\delta}^{(n)}) - E\nu_n(E_{(l+1)\delta}^{(n)} \setminus E_{l\delta}^{(n)})| > \frac{\eta}{6} k^{-1/2} (w_n(\gamma_1) \vee w_n(\gamma_2)) \right\} \\ & \leq \frac{k E\nu_n(E_{(l+1)\delta}^{(n)} \setminus E_{l\delta}^{(n)})}{(\eta/6)^2 k ((w_n(\gamma_1) \vee w_n(\gamma_2)))^2} \rightarrow 0 \end{aligned}$$

uniformly for all $l \in \{-\lceil M/\delta \rceil \dots \lfloor M/\delta \rfloor\}^6$.

It remains to prove that $Z_n(\theta) \rightarrow 0$ in probability for all $\theta \in [-M, M]^6$. This, however, follows similarly by Chebyshev's inequality, (D1) and the aforementioned analog to Corollary 4.5:

$$\begin{aligned}
& P\{|Z_n(\vartheta)| > \eta\} \\
&= P\left\{k|\nu_n(E_\theta^{(n)}) - E\nu_n(E_\theta^{(n)})| > \eta k^{1/2}(w_n(\gamma_1) \vee w_n(\gamma_2))\right\} \\
&\leq \frac{nP\{T_n^{\leftarrow}(X, Y) \in E_\theta^{(n)}\}}{\eta^2 k(w_n(\gamma_1) \vee w_n(\gamma_2))^2} \\
&= \frac{\nu(E_\theta^{(n)}) + O(A_0(n/k))}{\eta^2 (w_n(\gamma_1) \vee w_n(\gamma_2))^2} \\
&= \frac{\nu(S) + o(1)}{\eta^2 (w_n(\gamma_1) \vee w_n(\gamma_2))^2} \\
&\rightarrow 0.
\end{aligned}$$

□

Remark 4.7. Two remarks on this proof are in place. At first glance it seems peculiar that in the definition of $\tilde{H}_{\vartheta_i, \chi_i, \xi_i}^{(n, i)}$ both parameters $-\vartheta_i$ and χ_i take over the role of $k^{1/2}(\hat{\gamma}_i - \gamma_i)$ in the definition of \tilde{H} . This, however, is necessary to ensure the crucial monotonicity property of $\tilde{H}_{\vartheta_i, \chi_i, \xi_i}^{(n, i)}$ in the case $\gamma_i > 0$.

Secondly, we used the (slightly old-fashioned) classical approach to establish asymptotic equicontinuity instead of the often more elegant approach via bracketing numbers (see van der Vaart and Wellner (2000), Theorem 2.11.9), because the same approximation error of order $O(A_0(n/k))$ in (D1) always enters the upper bound on $E\nu_n(E_{(l+1)\delta}^{(n)} \setminus E_{l\delta}^{(n)})$, thus impeding the calculation of bracketing numbers for radii of smaller order.

Next we show that the terms *I* and *III* in decomposition (1.15) are negligible.

Lemma 4.8. *If the conditions of Theorem 2.1 are fulfilled, then*

$$\hat{p}_n - \frac{1}{e_n} \nu_n\left(\frac{d_n}{e_n} H_n(S_n^*)\right) = o_P(d_n^{-1} k^{-1/2}(w_n(\gamma_1) \vee w_n(\gamma_2))) \quad (4.24)$$

$$\frac{1}{e_n} (E\nu_n(B) - \nu(B))|_{B=(d_n/e_n)H_n(S_n^*)} = o_P(d_n^{-1} k^{-1/2}(w_n(\gamma_1) \vee w_n(\gamma_2))). \quad (4.25)$$

Proof. As $\hat{p}_n = \nu_n((d_n/e_n)H_n(S))/e_n$, the left-hand side of (4.24) is non-negative with expectation

$$\begin{aligned}
& \frac{n}{ke_n} P\left\{T_n^{\leftarrow}(X, Y) \in \frac{d_n}{e_n} H_n(S \setminus S_n^*)\right\} \\
&\leq \frac{n}{ke_n} P\left\{T_n^{\leftarrow}(X, Y) \in \frac{d_n}{e_n} H_n((0, u_n^*) \times [q(u_n^*), \infty) \cup [q^{\leftarrow}(v_n^*), \infty) \times [q(\infty), v_n^*)]\right\} \\
&= \frac{1}{d_n} \left(\nu(H_n((0, u_n^*) \times [q(u_n^*), \infty) \cup [q^{\leftarrow}(v_n^*), \infty) \times [q(\infty), v_n^*)]) + o(k^{-1/2}(w_n(\gamma_1) \vee w_n(\gamma_2)))\right)
\end{aligned}$$

where we have used (D1) and (S2). Now assertion (4.24) follows from Lemma 4.2 and the proof of Corollary 4.5.

Likewise, by conditions (D1), (S2) and $d_n \asymp e_n$, the left-hand side of (4.25) equals

$$\begin{aligned} \frac{1}{e_n} \left(\frac{n}{k} P\{T_n^{\leftarrow}(X, Y) \in B\} - \nu(B) \right) \Big|_{B=(d_n/e_n)H_n(S_n^*)} &= O_P(e_n^{-1} A_0(n/k)) \\ &= o_P(d_n^{-1} k^{-1/2} (w_n(\gamma_1) \vee w_n(\gamma_2))). \end{aligned}$$

□

Finally, we derive a bound on term VI in decomposition (1.15).

Lemma 4.9. *Under the assumptions of Theorem 2.1 one has*

$$\nu(d_n S) - p_n = o(d_n^{-1} k^{-1/2} (w_n(\gamma_1) \vee w_n(\gamma_2))).$$

Proof. With $\lambda_{n,i}, \tau_{n,i}$ as in Lemma 4.1, we define for $x \in [\lambda_{n,i}, \tau_{n,i}]$

$$H_{n,i}^*(x) := \left(1 + \gamma_i \frac{U_i(d_n x) - b_i(d_n)}{a_i(d_n)} \right)^{1/\gamma_i}.$$

According to de Haan and Ferreira (2006), Theorem 2.3.6 and 2.3.7 one can choose $a_i(t)$ as a multiple of t^{γ_i} and $b_i(t) = U_i(t) + O(a_i(t)A_i(t))$. Thus, for $\Delta_1(x)$ defined in the proof of Lemma 4.1

$$\begin{aligned} \frac{U_i(d_n x) - b_i(d_n)}{a_i(d_n)} &= \frac{a_i(n/k)}{a_i(d_n)} \left(\frac{U_i(d_n x) - b_i(n/k)}{a_i(n/k)} - \frac{b_i(d_n) - b_i(n/k)}{a_i(n/k)} \right) \\ &= \left(\frac{n}{k d_n} \right)^{\gamma_i} \left(\frac{(x d_n k/n)^{\gamma_i} - 1}{\gamma_i} + \Delta_1(x) + \frac{(d_n k/n)^{\gamma_i} - 1}{\gamma_i} + \Delta_1(1) \right) + O(A_i(d_n)) \\ &= \frac{x^{\gamma_i} - 1}{\gamma_i} + O\left(A_i(n/k) \left(\frac{d_n k}{n} \right)^{\rho_i + \varepsilon} (x^{\gamma_i + \rho_i + \varepsilon} + 1)\right) + o\left(A_i(n/k) \left(\frac{d_n k}{n} \right)^{\rho_i + \varepsilon}\right), \end{aligned}$$

where in the last step we have used (4.1), (4.18) and the Potter bound for the regularly varying function A_0 (de Haan and Ferreira (2006), Prop. B.1.9 5.). We conclude that

$$1 + \gamma_i \frac{U_i(d_n x) - b_i(d_n)}{a_i(d_n)} = x^{\gamma_i} \left(1 + O\left(A_i(n/k) \left(\frac{x d_n k}{n} \right)^{\rho_i + \varepsilon}\right) + O\left(A_i(n/k) \left(\frac{d_n k}{n} \right)^{\rho_i + \varepsilon} x^{-\gamma_i}\right) \right).$$

Check that the first remainder term is of smaller order than $k^{-1/2} w_n(\gamma_i)$ by condition (i) of Lemma 4.1. Moreover, for $\gamma_i > 0$, (4.19) and again condition (i) of Lemma 4.1 imply

$$A_i(n/k) \left(\frac{d_n k}{n} \right)^{\rho_i + \varepsilon} x^{-\gamma_i} = O\left(A_i(n/k) \left(\frac{d_n k}{n} \right)^{\rho_i + \varepsilon} k^{1/2} / \log c_n\right) \rightarrow 0,$$

while for $\gamma_i < 0$ this convergence follows from the conditions (i) and (iii) of Lemma 4.1, and for γ_i it is obvious from condition (i).

This shows that $H_{n,i}^*(x)$ is indeed well defined with

$$H_{n,i}^*(x) = x \left(1 + \begin{cases} o(k^{-1/2} \log c_n) + O\left(A_i(n/k) (d_n k/n)^{\rho_i + \varepsilon} x^{-\gamma_i}\right), & \gamma_i > 0 \\ o(k^{-1/2} (d_n k/n)^{-\gamma_i} (1 + x^{-\gamma_i})), & \gamma_i < 0 \\ o(k^{-1/2} \log^2 c_n), & \gamma_i = 0 \end{cases} \right)$$

uniformly for $x \in [\lambda_{n,i}, \tau_{n,i}]$. Notice that this representation is of similar type as the approximation derived in Lemma 4.1 with all leading terms equal to 0 (though in the case $\gamma_i > 0$ the second remainder term has a slightly different form). Therefore, we may proceed as before to conclude

$$\begin{aligned}\nu(H_n^*(S_n^*)) - \nu(S) &= o(k^{-1/2}(w_n(\gamma_1) \vee w_n(\gamma_2))) + \sum_{i=1}^2 O(A_i(n/k)(d_n k/n)^{\rho_i+\varepsilon}) 1_{\{\gamma_i > 0\}} \\ &= o(k^{-1/2}(w_n(\gamma_1) \vee w_n(\gamma_2))),\end{aligned}$$

where the last equality follows from Lemma 4.1 (i) (cf. Corollary 4.5).

To complete the proof, we must show that

$$p_n - \nu(d_n H_n^*(S_n^*)) = o(d_n^{-1} k^{-1/2}(w_n(\gamma_1) \vee w_n(\gamma_2))).$$

This, however, follows from assumption (D1) (with $t = d_n$) in a similar way as (4.24). \square

PROOF OF THEOREM 2.1. The assertion is a direct consequence of (1.15), Corollary 4.5 and of the Lemmas 4.8, 4.6 and 4.9. \square

PROOF OF COROLLARY 2.2. First note that, similarly as for \hat{p}_n , one obtains the representation $\hat{\nu}_n(\hat{S}_{n,2}^+) = \nu_n\left(\frac{d_n}{e_n} H_n^+(S)\right)$ with $H_n^+(x, y) := (H_{n,1}(x), H_{n,2}^+(y))$,

$$H_{n,2}^+(y) := \frac{e_n}{d_n} T_n^{\leftarrow} \circ \hat{T}_n \circ (\hat{T}_n^{(c^+)})^{\leftarrow} \circ U(d_n y)$$

and $c^+ := c_n^+ := (1 + \ell_n)n/(ke_n)$. Thus Lemma 4.1 (with e_n replaced by $e_n/(1 + \ell_n)$) yields the approximation

$$H_{n,2}^+(y) = (1 + \ell_n)y \left(1 + \begin{cases} -k^{-1/2} \log c_n(\Gamma_2/\gamma_2 + o_P(1)) + O_P(k^{-1/2}(y d_n/e_n)^{-\gamma_2}), & \gamma_2 > 0 \\ k^{-1/2}(d_n k/n)^{-\gamma_2}((\alpha_2/\gamma_2 - \beta_2 - \Gamma_2/\gamma_2^2 + o_P(1))y^{-\gamma_2} + o_P(1)), & \gamma_2 < 0 \\ -k^{-1/2} \log^2 c_n(\Gamma_2/2 + o_P(1)) + O_P(k^{-1/2} \log c_n \log y), & \gamma_2 = 0 \end{cases} \right)$$

Now the very same arguments as used in the analysis of \hat{p}_n show that

$$\hat{\nu}_n(\hat{S}_{n,2}^+) = \nu(S_{n,2}^+) + O_P(k^{-1/2}(w_n(\gamma_1) \vee w_n(\gamma_2))).$$

Together with an analogous approximation for $\hat{\nu}_n(\hat{S}_{n,2}^-)$ and our assumption on ℓ_n , we may conclude that

$$\begin{aligned}\frac{d_n}{e_n} \hat{I}_{n,2} &= \frac{d_n}{e_n} \frac{\nu(S_{n,2}^-) - \nu(S_{n,2}^+)}{2\ell_n} + o_P(1) \\ &= \int_{x_l}^{\infty} (2\ell_n)^{-1} \int_{(1-\ell_n)q(u)}^{(1+\ell_n)q(u)} \eta(u, v) dv du \\ &\rightarrow \int_{x_l}^{\infty} q(u) \eta(u, q(u)) du.\end{aligned}$$

In the last step we have used the fact that, on the range of integration, $\eta(u, v)$ is continuous and bounded by a multiple of $u^{-3} \vee (q(u))^{-3}$ (cf. (2.4)), so that the integrand of the outer integral can

easily be bounded by an integrable function and convergence follows by the dominated convergence theorem. \square

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